# Modern Real Analysis 

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## Preface

This text is an essentially self-contained treatment of material that is normally found in a first-year graduate course in real analysis. Although the presentation is based on a modern treatment of measure and integration, it has not lost sight of the fact that the theory of functions of one real variable is the core of the subject. It is assumed that the student has had a solid course in Advanced Calculus. Although the book's primary purpose is to serve as a graduate text, we hope that it will also serve as useful reference for the more experienced mathematician.

The book begins with a chapter on preliminaries and then proceeds with a chapter on the development of the real number system. This also includes an informal presentation of cardinal and ordinal numbers. The next chapter provides the basics of general topological and metric spaces. By the time this chapter has been concluded, the backgrounds of the students in a typical course will have been equalized and they will be prepared to pursue the main thrust of the book.

The text then proceeds to develop measure and integration theory in the next three chapters. Measure theory is introduced by first considering outer measures on an abstract space. The treatment here is abstract, yet short, simple, and basic. By focusing first on outer measures, the development underscores in a natural way the fundamental importance and significance of $\sigma$-algebras. Lebesgue measure, Lebesgue-Stieltjes measure, and Hausdorff measure are immediately developed as important, concrete examples of outer measures. Integration theory is presented by using countably simple functions, that is, functions that assume only a countable number of values. Conceptually they are no more difficult than simple functions, but their use leads to a more direct development. Important results such as the Radon-Nikodym theorem and Fubini's theorem have received treatments that avoid some of the usual technical difficulties.

A chapter on elementary functional analysis is followed by one on the Daniell integral and the Riesz Representation theorem. This introduces the student to a completely different approach to measure and integration theory. In order for the student to become more comfortable with this new framework, the linear functional
approach is further developed by including a short chapter on Schwartz Distributions. Along with introducing new ideas, this reinforces the student's previous encounter with measures as linear functionals. It also maintains connection with previous material by casting some old ideas in a new light. For example, $B V$ functions and absolutely continuous functions are characterized as functions whose distributional derivatives are measures and functions, respectively.

The introduction of Schwartz distributions invites a treatment of functions of several variables. Since absolutely continuous functions are so important in real analysis, it is natural to ask whether they have a counterpart among functions of several variables. In the last chapter, it is shown that this is the case by developing the class of functions whose partial derivatives (in the sense of distributions) are functions, thus providing a natural analog of absolutely continuous functions of a single variable. The analogy is strengthened by proving that these functions are absolutely continuous in each variable separately. These functions, called Sobolev functions, are of fundamental importance to many areas of research today. The chapter is concluded with a glimpse of both the power and the beauty of Distribution theory by providing a treatment of the Dirichlet Problem for Laplace's equation. This presentation is not difficult, but it does call upon many of the topics the student has learned throughout the text, thus providing a fitting end to the book.

We will use the following notation throughout. The symbol $\square$ denotes the end of a proof and $a:=b$ means $a=b$ by definition. All theorems, lemmas, corollaries, definitions, and remarks are numbered as $a . b$ where $a$ denotes the chapter number. Equation numbers are numbered in a similar way and appear as (a.b). Sections marked with * are not essential to the main development of the material and may be omitted.

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## CHAPTER 1

## Preliminaries

### 1.1. Sets


#### Abstract

This is the first of three sections devoted to basic definitions, notation, and terminology used throughout this book. We begin with an elementary and intuitive discussion of sets and deliberately avoid a rigorous treatment of "set theory" that would take us too far from our main purpose.


We shall assume that the notion of set is already known to the reader, at least in the intuitive sense. Roughly speaking, a set is any identifiable collection of objects called the elements or members of the set. Sets will usually be denoted by capital Roman letters such as $A, B, C, U, V, \ldots$, and if an object $x$ is an element of $A$, we will write $x \in A$. When $x$ is not an element of $A$ we write $x \notin A$. There are many ways in which the objects of a set may be identified. One way is to display all objects in the set. For example, $\left\{x_{1}, x_{2}, \ldots, x_{k}\right\}$ is the set consisting of the elements $x_{1}, x_{2}, \ldots, x_{k}$. In particular, $\{a, b\}$ is the set consisting of the elements $a$ and $b$. Note that $\{a, b\}$ and $\{b, a\}$ are the same set. A set consisting of a single element $x$ is denoted by $\{x\}$ and is called a singleton. Often it is possible to identify a set by describing properties that are possessed by its elements. That is, if $P(x)$ is a property possessed by an element $x$, then we write $\{x: P(x)\}$ to describe the set that consists of all objects $x$ for which the property $P(x)$ is true. Obviously, we have $A=\{x: x \in A\}$ and $\{x: x \neq x\}=\emptyset$, the empty set or null set.

The union of sets $A$ and $B$ is the set $\{x: x \in A$ or $x \in B\}$ and this is written as $A \cup B$. Similarly, if $\mathcal{A}$ is an arbitrary family of sets, the union of all sets in this family is

$$
\begin{equation*}
\{x: x \in A \text { for some } A \in \mathcal{A}\} \tag{1.1}
\end{equation*}
$$

and is denoted by

$$
\begin{equation*}
\bigcup_{A \in \mathcal{A}} A \quad \text { or as } \quad \bigcup\{A: A \in \mathcal{A}\} \tag{1.2}
\end{equation*}
$$

Sometimes a family of sets will be defined in terms of an indexing set $I$ and then we write

$$
\begin{equation*}
\left\{x: x \in A_{\alpha} \text { for some } \alpha \in I\right\}=\bigcup_{\alpha \in I} A_{\alpha} . \tag{1.3}
\end{equation*}
$$

If the index set $I$ is the set of positive integers, then we write (1.3) as

$$
\begin{equation*}
\bigcup_{i=1}^{\infty} A_{i} . \tag{1.4}
\end{equation*}
$$

The intersection of sets $A$ and $B$ is defined by $\{x: x \in A$ and $x \in B\}$ and is written as $A \cap B$. Similar to (1.1) and (1.2) we have

$$
\{x: x \in A \text { for all } A \in \mathcal{A}\}=\bigcap_{A \in \mathcal{A}} A=\bigcap\{A: A \in \mathcal{A}\}
$$

A family $\mathcal{A}$ of sets is said to be disjoint if $A_{1} \cap A_{2}=\emptyset$ for every pair $A_{1}$ and $A_{2}$ of distinct members of $\mathcal{A}$.

If every element of the set $A$ is also an element of $B$, then $A$ is called a subset of $B$ and this is written as $A \subset B$ or $B \supset A$. With this terminology, the possibility that $A=B$ is allowed. The set $A$ is called a proper subset of $B$ if $A \subset B$ and $A \neq B$.

The difference of two sets is

$$
A \backslash B=\{x: x \in A \text { and } x \notin B\}
$$

while the symmetric difference is

$$
A \Delta B=(A \backslash B) \cup(B \backslash A)
$$

In most discussions, a set $A$ will be a subset of some underlying space $X$ and in this context, we will speak of the complement of $A$ (relative to $X$ ) as the set $\{x: x \in X$ and $x \notin A\}$. This set is denoted by $\tilde{A}$ and this notation will be used if there is no doubt that complementation is taken with respect to $X$. In case of possible ambiguity, we write $X \backslash A$ instead of $\tilde{A}$. The following identities, known as de Morgan's laws, are very useful and easily verified:

$$
\begin{align*}
& \left(\bigcup_{\alpha \in I} A_{\alpha}\right)^{\sim}=\bigcap_{\alpha \in I} \widetilde{A_{\alpha}}  \tag{1.5}\\
& \left(\bigcap_{\alpha \in I} A_{\alpha}\right)^{\sim}=\bigcup_{\alpha \in I} \widetilde{A_{\alpha}} .
\end{align*}
$$

We shall denote the set of all subsets of $X$, called the power set of $X$, by $\mathcal{P}(X)$. Thus,

$$
\begin{equation*}
\mathcal{P}(X)=\{A: A \subset X\} \tag{1.6}
\end{equation*}
$$

The notions of limit superior (lim sup) and lim inferior (lim inf) are defined for sets as well as for sequences:

$$
\begin{align*}
\limsup _{i \rightarrow \infty} E_{i} & =\bigcap_{k=1}^{\infty} \bigcup_{i=k}^{\infty} E_{i}  \tag{1.7}\\
\liminf _{i \rightarrow \infty} E_{i} & =\bigcup_{k=1}^{\infty} \bigcap_{i=k}^{\infty} E_{i}
\end{align*}
$$

It is easily seen that

$$
\begin{align*}
\limsup _{i \rightarrow \infty} E_{i} & =\left\{x: x \in E_{i} \text { for infinitely many } i\right\}  \tag{1.8}\\
\liminf _{i \rightarrow \infty} E_{i} & =\left\{x: x \in E_{i} \text { for all but finitely many } i\right\} .
\end{align*}
$$

We use the following notation throughout:

$$
\begin{aligned}
& \emptyset=\text { the empty set } \\
& \mathbb{N}=\text { the set of positive integers, (not including zero), } \\
& \mathbb{Z}=\text { the set of integers, } \\
& \mathbb{Q}=\text { the set of rational numbers }, \\
& \mathbb{R}=\text { the set of real numbers. }
\end{aligned}
$$

We assume the reader has knowledge of the sets $\mathbb{N}, \mathbb{Z}$, and $\mathbb{Q}$, while $\mathbb{R}$ will be carefully constructed in Section 2.1.

## Exercises for Section 1.1

1. Two sets are identical if and only if they have the same members. That is, $A=B$ if and only if for each element $x, x \in A$ when and only when $x \in B$. Prove $A=B$ if and only if $A \subset B$ and $B \subset A$.
2. Prove that $A \subset B$ if and only if $A=A \cup B$.
3. Prove de Morgan's laws, (1.5).
4. Let $E_{i}, i=1,2, \ldots$, be a family of sets. Use definitions (1.7) to prove

$$
\liminf _{i \rightarrow \infty} E_{i} \subset \limsup _{i \rightarrow \infty} E_{i}
$$

### 1.2. Functions

In this section an informal discussion of relations and functions is given, a subject that is encountered in several forms in elementary analysis. In this development, we adopt the notion that a relation or function is indistinguishable from its graph.

If $X$ and $Y$ are sets, the Cartesian product of $X$ and $Y$ is

$$
\begin{equation*}
X \times Y=\{\text { all ordered pairs }(x, y): x \in X, y \in Y\} \tag{1.9}
\end{equation*}
$$

The ordered pair $(x, y)$ is thus to be distinguished from $(y, x)$. We will discuss the Cartesian product of an arbitrary family of sets later in this section.

A relation from $X$ to $Y$ is a subset of $X \times Y$. If $f$ is a relation, then the domain and range of $f$ are

$$
\begin{aligned}
\operatorname{dom} f & =X \cap\{x:(x, y) \in f \text { for some } y \in Y\} \\
\operatorname{rng} f & =Y \cap\{y:(x, y) \in f \text { for some } x \in X\} .
\end{aligned}
$$

Frequently symbols such as $\sim$ or $\leq$ are used to designate a relation. In these cases the notation $x \sim y$ or $x \leq y$ will mean that the element $(x, y)$ is a member of the relation $\sim$ or $\leq$, respectively.

A relation $f$ is said to be single-valued if $y=z$ whenever $(x, y)$ and $(x, z) \in$ $f$. A single-valued relation is called a function. The terms mapping, map, transformation are frequently used interchangeably with function, although the term function is usually reserved for the case when the range of $f$ is a subset of $\mathbb{R}$. If $f$ is a mapping and $(x, y) \in f$, we let $f(x)$ denote $y$. We call $f(x)$ the image of $x$ under $f$. We will also use the notation $x \mapsto f(x)$, which indicates that $x$ is mapped to $f(x)$ by $f$. If $A \subset X$, then the image of $A$ under $f$ is

$$
\begin{equation*}
f(A)=\{y: y=f(x), \text { for some } x \in \operatorname{dom} f \cap A\} \tag{1.10}
\end{equation*}
$$

Also, the inverse image of $B$ under $f$ is

$$
\begin{equation*}
f^{-1}(B)=\{x: x \in \operatorname{dom} f, f(x) \in B\} \tag{1.11}
\end{equation*}
$$

In case the set $B$ consists of a single point $y$, or in other words $B=\{y\}$, we will simply write $f^{-1}\{y\}$ instead of the full notation $f^{-1}(\{y\})$. If $A \subset X$ and $f$ a mapping with $\operatorname{dom} f \subset X$, then the restriction of $f$ to $A$, denoted by $f L A$, is defined by $f L A(x)=f(x)$ for all $x \in A \cap \operatorname{dom} f$.

If $f$ is a mapping from $X$ to $Y$ and $g$ a mapping from $Y$ to $Z$, then the composition of $g$ with $f$ is a mapping from $X$ to $Z$ defined by

$$
\begin{equation*}
g \circ f=\{(x, z):(x, y) \in f \text { and }(y, z) \in g \text { for some } y \in Y\} \tag{1.12}
\end{equation*}
$$

If $f$ is a mapping such that $\operatorname{dom} f=X$ and $\operatorname{rng} f \subset Y$, then we write $f: X \rightarrow Y$. The mapping $f$ is called an injection or is said to be univalent if $f(x) \neq f\left(x^{\prime}\right)$ whenever $x, x^{\prime} \in \operatorname{dom} f$ with $x \neq x^{\prime}$. The mapping $f$ is called a surjection or onto $Y$ if for each $y \in Y$, there exists $x \in X$ such that $f(x)=y$; in other words, $f$ is a surjection if $f(X)=Y$. Finally, we say that $f$ is a bijection if $f$ is both an injection and a surjection. A bijection $f: X \rightarrow Y$ is also called a one-to-one correspondence between $X$ and $Y$.

There is one relation that is particularly important and is so often encountered that it requires a separate definition.
1.1. Definition. If $X$ is a set, an equivalence relation on $X$ (often denoted by $\sim$ ) is a relation characterized by the following conditions:
(i) $x \sim x$ for every $x \in X$
(reflexive)
(ii) if $x \sim y$, then $y \sim x$,
(symmetric)
(iii) if $x \sim y$ and $y \sim z$, then $x \sim z$.
(transitive)

Given an equivalence relation $\sim$ on $X$, a subset $A$ of $X$ is called an equivalence class if and only if there is an element $x \in A$ such that $A$ consists precisely of those elements $y$ such that $x \sim y$. One can easily verify that dsitinct equivalence classes are disjoint and that $X$ can be expressed as the union of equivalence classes.

A sequence in a space $X$ is a mapping $f: \mathbb{N} \rightarrow X$. It is convenient to denote a sequence $f$ as a list. Thus, if $f(k)=x_{k}$, we speak of the sequence $\left\{x_{k}\right\}_{k=1}^{\infty}$ or simply $\left\{x_{k}\right\}$. A subsequence is obtained by discarding some elements of the original sequence and ordering the elements that remain in the usual way. Formally, we say that $x_{k_{1}}, x_{k_{2}}, x_{k_{3}}, \ldots$, is a subsequence of $x_{1}, x_{2}, x_{3}, \ldots$, if there is a mapping $g: \mathbb{N} \rightarrow \mathbb{N}$ such that for each $i \in \mathbb{N}, x_{k_{i}}=x_{g(i)}$ and if $g(i)<g(j)$ whenever $i<j$.

Our final topic in this section is the Cartesian product of a family of sets. Let $\mathcal{X}$ be a family of sets $X_{\alpha}$ indexed by a set $I$. The Cartesian product of $\mathcal{X}$ is denoted by

$$
\prod_{\alpha \in I} X_{\alpha}
$$

and is defined as the set of all mappings

$$
x: I \rightarrow \bigcup X_{\alpha}
$$

with the property that

$$
\begin{equation*}
x(\alpha) \in X_{\alpha} \tag{1.13}
\end{equation*}
$$

for each $\alpha \in I$. Each mapping $x$ is called a choice mapping for the family $\mathcal{X}$. Also, we call $x(\alpha)$ the $\alpha^{\text {th }}$ coordinate of $x$. This terminology is perhaps easier to understand if we consider the case where $I=\{1,2, \ldots, n\}$. As in the preceding paragraph, it is useful to denote the choice mapping $x$ as a list $\{x(1), x(2), \ldots, x(n)\}$, and even more useful if we write $x(i)=x_{i}$. The mapping $x$ is thus identified with the ordered $n$-tuple $\left(x_{1}, x_{2}, \ldots, x_{n}\right)$. Here, the word "ordered" is crucial because an $n$-tuple consisting of the same elements but in a different order produces a different mapping $x$. Consequently, the Cartesian product becomes the set of all ordered
$n$-tuples:

$$
\begin{equation*}
\prod_{i=1}^{n} X_{i}=\left\{\left(x_{1}, x_{2}, \ldots, x_{n}\right): x_{i} \in X_{i}, i=1,2, \ldots, n\right\} \tag{1.14}
\end{equation*}
$$

In the special case where $X_{i}=\mathbb{R}, i=1,2, \ldots, n$, an element of the Cartesian product is a mapping that can be identified with an ordered $n$-tuple of real numbers. We denote the set of all ordered $n$-tuples (also referred to as vectors) by

$$
\mathbb{R}^{n}=\left\{\left(x_{1}, x_{2}, \ldots, x_{n}\right): x_{i} \in \mathbb{R}, i=1,2, \ldots, n\right\}
$$

$\mathbb{R}^{n}$ is called Euclidean $n$-space. The norm of a vector $x$ is defined as

$$
\begin{equation*}
|x|=\sqrt{x_{1}^{2}+x_{2}^{2}+\cdots+x_{n}^{2}} ; \tag{1.15}
\end{equation*}
$$

the distance between two vectors $x$ and $y$ is $|x-y|$. As we mentioned earlier in this section, the Cartesian product of two sets $X_{1}$ and $X_{2}$ is denoted by $X_{1} \times X_{2}$.
1.2. Remark. A fundamental issue that we have not addressed is whether the Cartesian product of an arbitrary family of sets is nonempty. This involves concepts from set theory and is the subject of the next section.

## Exercises for Section 1.2

1. Prove that $f \circ(g \circ h)=(f \circ g) \circ h$ for mappings $f, g$, and $h$.
2. Prove that $(f \circ g)^{-1}(A)=g^{-1}\left[f^{-1}(A)\right]$ for mappings $f$ and $g$ and an arbitrary set $A$.
3. Prove: If $f: X \rightarrow Y$ is a mapping and $A \subset B \subset X$, then $f(A) \subset f(B) \subset Y$. Also, prove that if $E \subset F \subset Y$, then $f^{-1}(E) \subset f^{-1}(F) \subset X$.
4. Prove: If $\mathcal{A} \subset \mathcal{P}(X)$, then

$$
f\left(\bigcup_{A \in \mathcal{A}} A\right)=\bigcup_{A \in \mathcal{A}} f(A) \text { and } f\left(\bigcap_{A \in \mathcal{A}} A\right) \subset \bigcap_{A \in \mathcal{A}} f(A)
$$

and

$$
f^{-1}\left(\bigcup_{A \in \mathcal{A}} A\right)=\bigcup_{A \in \mathcal{A}} f^{-1}(A) \text { and } f^{-1}\left(\bigcap_{A \in \mathcal{A}} A\right)=\bigcap_{A \in \mathcal{A}} f^{-1}(A) .
$$

Give an example that shows the above inclusion cannot be replaced by equality.
5. Consider a nonempty set $X$ and its power set $\mathcal{P}(X)$. For each $x \in X$, let $B_{x}=\{0,1\}$ and consider the Cartesian product $\prod_{x \in X} B_{x}$. Exhibit a natural one-to-one correspondence between $\mathcal{P}(X)$ and $\prod_{x \in X} B_{x}$.
6. Let $X \xrightarrow{f} Y$ be an arbitrary mapping and suppose there is a mapping $Y \xrightarrow{g} X$ such that $f \circ g(y)=y$ for all $y \in Y$ and that $g \circ f(x)=x$ for all $x \in X$. Prove that $f$ is one-to-one from $X$ onto $Y$ and that $g=f^{-1}$.
7. Show that $A \times(B \cup C)=(A \times B) \cup(A \times C)$. Also, show that in general $A \cup(B \times C) \neq(A \cup B) \times(A \cup C)$.

### 1.3. Set Theory

The material discussed in the previous two sections is based on tools found in elementary set theory. However, in more advanced areas of mathematics this material is not sufficient to discuss or even formulate some of the concepts that are needed. An example of this occurred in the previous section during the discussion of the Cartesian product of an arbitrary family of sets. Indeed, the Cartesian product of families of sets requires the notion of a choice mapping whose existence is not obvious. Here, we give a brief review of the Axiom of Choice and some of its logical equivalences.

A fundamental question that arises in the definition of the Cartesian product of an arbitrary family of sets is the existence of choice mappings. This is an example of a question that cannot be answered within the context of elementary set theory. In the beginning of the $20^{\text {th }}$ century, Ernst Zermelo formulated an axiom of set theory called the Axiom of Choice, which asserts that the Cartesian product of an arbitrary family of nonempty sets exists and is nonempty. The formal statement is as follows.
1.3. The Axiom of Choice. If $X_{\alpha}$ is a nonempty set for each element $\alpha$ of an index set $I$, then

$$
\prod_{\alpha \in I} X_{\alpha}
$$

is nonempty.
1.4. Proposition. The following statement is equivalent to the Axiom of Choice: If $\left\{X_{\alpha}\right\}_{\alpha \in A}$ is a disjoint family of nonempty sets, then there is a set $S \subset \cup_{\alpha \in A} X_{\alpha}$ such that $S \cap X_{\alpha}$ consists of precisely one element for every $\alpha \in A$.

Proof. The Axiom of Choice states that there exists $f: A \rightarrow \cup_{\alpha \in A} X_{\alpha}$ such that $f(\alpha) \in X_{\alpha}$ for each $\alpha \in A$. The set $S:=f(A)$ satisfies the conclusion of the statement. Conversely, if such a set $S$ exists, then the mapping $A \xrightarrow{f} \cup_{\alpha \in A} X_{\alpha}$ defined by assigning the point $S \cap X_{\alpha}$ the value of $f(\alpha)$ implies the validity of the Axiom of Choice.
1.5. Definition. Given a set $S$ and a relation $\leq$ on S , we say that $\leq$ is a partial ordering if the following three conditions are satisfied:
(i) $x \leq x$ for every $x \in S$

> (reflexive)
(ii) if $x \leq y$ and $y \leq x$, then $x=y$,
(antisymmetric)
(iii) if $x \leq y$ and $y \leq z$, then $x \leq z$. (transitive)
If, in addition,
(iv) either $x \leq y$ or $y \leq x$, for all $x, y \in S$,
(trichotomy)
then $\leq$ is called a linear or total ordering.
For example, $\mathbb{Z}$ is linearly ordered with its usual ordering, whereas the family of all subsets of a given set $X$ is partially ordered (but not linearly ordered) by $\subset$. If a set $X$ is endowed with a linear ordering, then each subset $A$ of $X$ inherits the ordering of $X$. That is, the restriction to $A$ of the linear ordering on $X$ induces a linear ordering on $A$. The following two statements are known to be equivalent to the Axiom of Choice.
1.6. Hausdorff Maximal Principle. Every partially ordered set has a maximal linearly ordered subset.
1.7. Zorn's Lemma. If $X$ is a partially ordered set with the property that each linearly ordered subset has an upper bound, then $X$ has a maximal element. In particular, this implies that if $\mathcal{E}$ is a family of sets (or a collection of families of sets) and if $\{\cup F: F \in \mathcal{F}\} \in \mathcal{E}$ for any subfamily $\mathcal{F}$ of $\mathcal{E}$ with the property that

$$
F \subset G \quad \text { or } \quad G \subset F \quad \text { whenever } \quad F, G \in \mathcal{F},
$$

then there exists $E \in \mathcal{E}$, which is maximal in the sense that it is not a subset of any other member of $\mathcal{E}$.

In the following, we will consider other formulations of the Axiom of Choice. This will require the notion of a linear ordering on a set.

A non-empty set $X$ endowed with a linear order is said to be well ordered if each subset of $X$ has a first element with respect to its induced linear order. Thus, the integers, $\mathbb{Z}$, with the usual ordering is not a well-ordered set, whereas the set $\mathbb{N}$ is well ordered. However, it is possible to define a linear ordering on $\mathbb{Z}$ that produces a well ordering. In fact, it is possible to do this for an arbitrary set if we assume the validity of the Axiom of Choice. This is stated formally in the Well-Ordering Theorem.
1.8. Theorem (The Well-Ordering Theorem). Every set can be well ordered. That is, if $A$ is an arbitrary set, then there exists a linear ordering of $A$ with the property that each non-empty subset of $A$ has a first element.

Cantor put forward the continuum hypothesis in 1878, conjecturing that every infinite subset of the continuum is either countable (i.e., can be put in 1-1 correspondence with the natural numbers) or has the cardinality of the continuum (i.e., can be put in 1-1 correspondence with the real numbers). The importance of this was seen by Hilbert who made the continuum hypothesis the first in the list of problems which he proposed in his Paris lecture of 1900. Hilbert saw this as one of the most fundamental questions which mathematicians should attack in the 1900s and he went further in proposing a method to attack the conjecture. He suggested that first one should try to prove another of Cantor's conjectures, namely that any set can be well ordered.

Zermelo began to work on the problems of set theory by pursuing, in particular, Hilbert's idea of resolving the problem of the continuum hypothesis. In 1902 Zermelo published his first work on set theory which was on the addition of transfinite cardinals. Two years later, in 1904, he succeeded in taking the first step suggested by Hilbert towards the continuum hypothesis when he proved that every set can be well ordered. This result brought fame to Zermelo and also earned him a quick promotion; in December 1905, he was appointed as professor in Göttingen.

The axiom of choice is the basis for Zermelo's proof that every set can be well ordered; in fact the axiom of choice is equivalent to the well-ordering property so we now know that this axiom must be used. His proof of the well-ordering property used the axiom of choice to construct sets by transfinite induction. Although Zermelo certainly gained fame for his proof of the well ordering property, set theory at this time was in the rather unusual position that many mathematicians rejected the type of proofs that Zermelo had discovered. There were strong feelings as to whether such non-constructive parts of mathematics were legitimate areas for study and Zermelo's ideas were certainly not accepted by quite a number of mathematicians.

The fundamental discoveries of K. Gödel [31] and P. J. Cohen [15], [17] shook the foundations of mathematics with results that placed the axiom of choice in a very interesting position. Their work shows that the Axiom of Choice, in fact, is a new principle in set theory because it can neither be proved nor disproved from the usual Zermelo-Fraenkel axioms of set theory. Indeed, Gödel showed, in 1940, that the Axiom of Choice cannot be disproved using the other axioms of set theory and then in 1963, Paul Cohen proved that the Axiom of Choice is independent of the other axioms of set theory. The importance of the Axiom of Choice will readily be seen throughout the following development, as we appeal to it in a variety of
contexts.

## Exercises for Section 1.3

1. Use a one-to-one correspondence between $\mathbb{Z}$ and $\mathbb{N}$ to exhibit a linear ordering of $\mathbb{N}$ that is not a well ordering.
2. Use the natural partial ordering of $\mathcal{P}(\{1,2,3\})$ to exhibit a partial ordering of $\mathbb{N}$ that is not a linear ordering.
3. For $(a, b),(c, d) \in \mathbb{N} \times \mathbb{N}$, define $(a, b) \leq(c, d)$ if either $a<c$ or $a=c$ and $b \leq d$. With this relation, prove that $\mathbb{N} \times \mathbb{N}$ is a well-ordered set.
4. Let $P$ denote the space of all polynomials defined on $\mathbb{R}$. For $p_{1}, p_{2} \in P$, define $p_{1} \leq p_{2}$ if there exists $x_{0}$ such that $p_{1}(x) \leq p_{2}(x)$ for all $x \geq x_{0}$. Is $\leq$ a linear ordering? Is $P$ well ordered?
5. Let $C$ denote the space of all continuous functions on $[0,1]$. For $f_{1}, f_{2} \in C$, define $f_{1} \leq f_{2}$ if $f_{1}(x) \leq f_{2}(x)$ for all $x \in[0,1]$. Is $\leq$ a linear ordering? Is $C$ well ordered?
6. Prove that the following assertion is equivalent to the Axiom of Choice: If $A$ and $B$ are nonempty sets and $f: A \rightarrow B$ is a surjection (that is, $f(A)=B$ ), then there exists a function $g: B \rightarrow A$ such that $g(y) \in f^{-1}(y)$ for each $y \in B$.
7. Use the following outline to prove that for any two sets $A$ and $B$, either card $A \leq$ $\operatorname{card} B$ or card $B \leq \operatorname{card} A$ : Let $\mathcal{F}$ denote the family of all injections from subsets of $A$ into $B$. Since $\mathcal{F}$ can be considered as a family of subsets of $A \times B$, it can be partially ordered by inclusion. Thus, we can apply Zorn's lemma to conclude that $\mathcal{F}$ has a maximal element, say $f$. If $a \in A \backslash \operatorname{dom} f$ and $b \in B \backslash f(A)$, then extend $f$ to $A \cup\{a\}$ by defining $f(a)=b$. Then $f$ remains an injection and thus contradicts maximality. Hence, either $\operatorname{dom} f=A$ in which case card $A \leq \operatorname{card} B$ or $B=\operatorname{rng} f$ in which case $f^{-1}$ is an injection from $B$ into $A$, which would imply $\operatorname{card} B \leq \operatorname{card} A$.
8. Complete the details of the following proposition: If card $A \leq \operatorname{card} B$ and $\operatorname{card} B \leq \operatorname{card} A$, then $\operatorname{card} A=\operatorname{card} B$.

Let $f: A \rightarrow B$ and $g: B \rightarrow A$ be injections. If $a \in A \cap \operatorname{rng} g$, we have $g^{-1}(a) \in B$. If $g^{-1}(a) \in \operatorname{rng} f$, we have $f^{-1}\left(g^{-1}(a)\right) \in A$. Continue this process as far as possible. There are three possibilities: either the process continues indefinitely, or it terminates with an element of $A \backslash \operatorname{rng} g$ (possibly with $a$ itself) or it terminates with an element of $B \backslash \operatorname{rng} f$. These three cases determine disjoint sets $A_{\infty}, A_{A}$ and $A_{B}$ whose union is $A$. In a similar manner, $B$ can be decomposed into $B_{\infty}, B_{B}$ and $B_{A}$. Now $f$ maps $A_{\infty}$ onto $B_{\infty}$ and $A_{A}$ onto $B_{A}$
and $g$ maps $B_{B}$ onto $A_{B}$. If we define $h: A \rightarrow B$ by $h(a)=f(a)$ if $a \in A_{\infty} \cup A_{A}$ and $h(a)=g^{-1}(a)$ if $a \in A_{B}$, we find that $h$ is injective.

## CHAPTER 2

## Real, Cardinal and Ordinal Numbers

### 2.1. The Real Numbers

A brief development of the construction of the Real Numbers is given in terms of equivalence classes of Cauchy sequences of rational numbers. This construction is based on the assumption that properties of the rational numbers, including the integers, are known.

In our development of the real number system, we shall assume that properties of the natural numbers, integers, and rational numbers are known. In order to agree on what the properties are, we summarize some of the more basic ones. Recall that the natural numbers are designated as

$$
\mathbb{N}:=\{1,2, \ldots, k, \ldots\}
$$

They form a well-ordered set when endowed with the usual ordering. The ordering on $\mathbb{N}$ satisfies the following properties:
(i) $x \leq x$ for every $x \in S$.
(ii) if $x \leq y$ and $y \leq x$, then $x=y$.
(iii) if $x \leq y$ and $y \leq z$, then $x \leq z$.
(iv) for all $x, y \in S$, either $x \leq y$ or $y \leq x$.

The four conditions above define a linear ordering on $S$, a topic that was introduced in Section 1.3 and will be discussed in greater detail in Section 2.3. The linear order $\leq$ of $\mathbb{N}$ is compatible with the addition and multiplication operations in $\mathbb{N}$. Furthermore, the following three conditions are satisfied:
(i) Every nonempty subset of $\mathbb{N}$ has a first element; i.e., if $\emptyset \neq S \subset \mathbb{N}$, there is an element $x \in S$ such that $x \leq y$ for any element $y \in S$. In particular, the set $\mathbb{N}$ itself has a first element that is unique, in view of (ii) above, and is denoted by the symbol 1 ,
(ii) Every element of $\mathbb{N}$, except the first, has an immediate predecessor. That is, if $x \in \mathbb{N}$ and $x \neq 1$, then there exists $y \in \mathbb{N}$ with the property that $y \leq x$ and $z \leq y$ whenever $z \leq x$.
(iii) $\mathbb{N}$ has no greatest element; i.e., for every $x \in \mathbb{N}$, there exists $y \in \mathbb{N}$ such that $x \neq y$ and $x \leq y$.

The reader can easily show that (i) and (iii) imply that each element of $\mathbb{N}$ has an immediate successor, i.e., that for each $x \in \mathbb{N}$, there exists $y \in \mathbb{N}$ such that $x<y$ and that if $x<z$ for some $z \in \mathbb{N}$ where $y \neq z$, then $y<z$. The immediate successor of $x, y$, will be denoted by $x^{\prime}$. A nonempty set $S \subset \mathbb{N}$ is said to be finite if $S$ has a greatest element.

From the structure established above follows an extremely important result, the so-called principle of mathematical induction, which we now prove.
2.1. Theorem. Suppose $S \subset \mathbb{N}$ is a set with the property that $1 \in S$ and that $x \in S$ implies $x^{\prime} \in S$. Then $S=\mathbb{N}$.

Proof. Suppose $S$ is a proper subset of $\mathbb{N}$ that satisfies the hypotheses of the theorem. Then $\mathbb{N} \backslash S$ is a nonempty set and therefore by (i) above, has a first element $x$. Note that $x \neq 1$ since $1 \in S$. From (ii) we see that $x$ has an immediate predecessor, $y$. As $y \in S$, we have $y^{\prime} \in S$. Since $x=y^{\prime}$, we have $x \in S$, contradicting the choice of $x$ as the first element of $\mathbb{N} \backslash S$.

Also, we have $x \in S$ since $x=y^{\prime}$. By definition, $x$ is the first element of $\mathbb{N}-S$, thus producing a contradiction. Hence, $S=\mathbb{N}$.

The rational numbers $\mathbb{Q}$ may be constructed in a formal way from the natural numbers. This is accomplished by first defining the integers, both negative and positive, so that subtraction can be performed. Then the rationals are defined using the properties of the integers. We will not go into this construction but instead leave it to the reader to consult another source for this development. We list below the basic properties of the rational numbers.

The rational numbers are endowed with the operations of addition and multiplication that satisfy the following conditions:
(i) For every $r, s \in \mathbb{Q}, r+s \in \mathbb{Q}$, and $r s \in \mathbb{Q}$.
(ii) Both operations are commutative and associative, i.e., $r+s=s+r, r s=$ $s r,(r+s)+t=r+(s+t)$, and $(r s) t=r(s t)$.
(iii) The operations of addition and multiplication have identity elements 0 and 1 respectively, i.e., for each $r \in \mathbb{Q}$, we have

$$
0+r=r \quad \text { and } \quad 1 \cdot r=r
$$

(iv) The distributive law is valid:

$$
r(s+t)=r s+r t
$$

whenever $r, s$, and $t$ are elements of $\mathbb{Q}$.
(v) The equation $r+x=s$ has a solution for every $r, s \in \mathbb{Q}$. The solution is denoted by $s-r$.
(vi) The equation $r x=s$ has a solution for every $r, s \in \mathbb{Q}$ with $r \neq 0$. This solution is denoted by $s / r$. Any set containing at least two elements and satisfying the six conditions above is called a field; in particular, the rational numbers form a field. The set $\mathbb{Q}$ can also be endowed with a linear ordering. The order relation is related to the operations of addition and multiplication as follows:
(vii) If $r \geq s$, then for every $t \in \mathbb{Q}, r+t \geq s+t$.
(viii) $0<1$.
(ix) If $r \geq s$ and $t \geq 0$, then $r t \geq s t$.

The rational numbers thus provides an example of an ordered field. The proof of the following is elementary and is left to the reader, see Exercise 6 at the end of this section.
2.2. Theorem. Every ordered field $F$ contains an isomorphic image of $\mathbb{Q}$ and the isomorphism can be taken as order preserving.

In view of this result, we may view $\mathbb{Q}$ as a subset of $F$. Consequently, the following definition is meaningful.
2.3. Definition. An ordered field $F$ is called an Archimedean ordered field, if for each $a \in F$ and each positive $b \in \mathbb{Q}$, there exists a positive integer $n$ such that $n b>a$. Intuitively, this means that no matter how large $a$ is and how small $b$, successive repetitions of $b$ will eventually exceed $a$.

Although the rational numbers form a rich algebraic system, they are inadequate for the purposes of analysis because they are, in a sense, incomplete. For example, a negative rational number does not have a rational square root, and not every positive rational number has a rational square root. We now proceed to construct the real numbers assuming knowledge of the integers and rational numbers. This is basically an assumption concerning the algebraic structure of the real numbers.

The linear order structure of the field permits us to define the notion of the absolute value of an element of the field. That is, the absolute value of $x$ is defined by

$$
|x|=\left\{\begin{aligned}
x & \text { if } x \geq 0 \\
-x & \text { if } x<0
\end{aligned}\right.
$$

We will freely use properties of the absolute value such as the triangle inequality in our development.

The following two definitions are undoubtedly well known to the reader; we state them only to emphasize that at this stage of the development, we assume knowledge of only the rational numbers.
2.4. Definition. A sequence of rational numbers $\left\{r_{i}\right\}$ is Cauchy if and only if for each rational $\varepsilon>0$, there exists a positive integer $N(\varepsilon)$ such that $\left|r_{i}-r_{k}\right|<\varepsilon$ whenever $i, k \geq N(\varepsilon)$.
2.5. Definition. A rational number $r$ is said to be the limit of a sequence of rational numbers $\left\{r_{i}\right\}$ if and only if for each rational $\varepsilon>0$, there exists a positive integer $N(\varepsilon)$ such that

$$
\left|r_{i}-r\right|<\varepsilon
$$

for $i \geq N(\varepsilon)$. This is written as

$$
\lim _{i \rightarrow \infty} r_{i}=r
$$

and we say that $\left\{r_{i}\right\}$ converges to $r$.
We leave the proof of the following proposition to the reader.
2.6. Proposition. A sequence of rational numbers that converges to a rational number is Cauchy.
2.7. Proposition. A Cauchy sequence of rational numbers, $\left\{r_{i}\right\}$, is bounded. That is, there exists a rational number $M$ such that $\left|r_{i}\right| \leq M$ for $i=1,2, \ldots$.

Proof. Choose $\varepsilon=1$. Since the sequence $\left\{r_{i}\right\}$ is Cauchy, there exists a positive integer $N$ such that

$$
\left|r_{i}-r_{j}\right|<1 \quad \text { whenever } \quad i, j \geq N
$$

In particular, $\left|r_{i}-r_{N}\right|<1$ whenever $i \geq N$. By the triangle inequality, $\left|r_{i}\right|-\left|r_{N}\right| \leq$ $\left|r_{i}-r_{N}\right|$ and therefore,

$$
\left|r_{i}\right|<\left|r_{N}\right|+1 \quad \text { for all } \quad i \geq N .
$$

If we define

$$
M=\operatorname{Max}\left\{\left|r_{1}\right|,\left|r_{2}\right|, \ldots,\left|r_{N-1}\right|,\left|r_{N}\right|+1\right\}
$$

then $\left|r_{i}\right| \leq M$ for all $i \geq 1$.
The reader can easily provide a proof of the following.
2.8. Proposition. Every Cauchy sequence of rational numbers has at most one limit.

The fact that some Cauchy sequences in $\mathbb{Q}$ do not have a limit (in $\mathbb{Q}$ ) is what makes $\mathbb{Q}$ incomplete. We will construct the completion by means of equivalence classes of Cauchy sequences.
2.9. Definition. Two Cauchy sequences of rational numbers $\left\{r_{i}\right\}$ and $\left\{s_{i}\right\}$ are said to be equivalent if and only if

$$
\lim _{i \rightarrow \infty}\left(r_{i}-s_{i}\right)=0
$$

We write $\left\{r_{i}\right\} \sim\left\{s_{i}\right\}$ when $\left\{r_{i}\right\}$ and $\left\{s_{i}\right\}$ are equivalent. It is easy to show that this, in fact, is an equivalence relation. That is,
(i) $\left\{r_{i}\right\} \sim\left\{r_{i}\right\}$, (reflexivity)
(ii) $\left\{r_{i}\right\} \sim\left\{s_{i}\right\}$ if and only if $\left\{s_{i}\right\} \sim\left\{r_{i}\right\}$,
(symmetry)
(iii) if $\left\{r_{i}\right\} \sim\left\{s_{i}\right\}$ and $\left\{s_{i}\right\} \sim\left\{t_{i}\right\}$, then $\left\{r_{i}\right\} \sim\left\{t_{i}\right\}$. (transitivity)
The set of all Cauchy sequences of rational numbers equivalent to a fixed Cauchy sequence is called an equivalence class of Cauchy sequences. The fact that we are dealing with an equivalence relation implies that the set of all Cauchy sequences of rational numbers is partitioned into mutually disjoint equivalence classes. For each rational number $r$, the sequence each of whose values is $r$ (i.e., the constant sequence) will be denoted by $\bar{r}$. Hence, $\overline{0}$ is the constant sequence whose values are 0 . This brings us to the definition of a real number.
2.10. Definition. An equivalence class of Cauchy sequences of rational numbers is termed a real number. In this section, we will usually denote real numbers by $\rho, \sigma$, etc. With this convention, a real number $\rho$ designates an equivalence class of Cauchy sequences, and if this equivalence class contains the sequence $\left\{r_{i}\right\}$, we will write

$$
\rho=\overline{\left\{r_{i}\right\}}
$$

and say that $\rho$ is represented by $\left\{r_{i}\right\}$. Note that $\{1 / i\}_{i=1}^{\infty} \sim \overline{0}$ and that every $\rho$ has a representative $\left\{r_{i}\right\}_{i=1}^{\infty}$ with $r_{i} \neq 0$ for every $i$.

In order to define the sum and product of real numbers, we invoke the corresponding operations on Cauchy sequences of rational numbers. This will require the next two elementary propositions whose proofs are left to the reader.
2.11. Proposition. If $\left\{r_{i}\right\}$ and $\left\{s_{i}\right\}$ are Cauchy sequences of rational numbers, then $\left\{r_{i} \pm s_{i}\right\}$ and $\left\{r_{i} \cdot s_{i}\right\}$ are Cauchy sequences. The sequence $\left\{r_{i} / s_{i}\right\}$ is also Cauchy provided $s_{i} \neq 0$ for every $i$ and $\left\{s_{i}\right\}_{i=1}^{\infty} \nsim \overline{0}$.
2.12. Proposition. If $\left\{r_{i}\right\} \sim\left\{r_{i}^{\prime}\right\}$ and $\left\{s_{i}\right\} \sim\left\{s_{i}^{\prime}\right\}$, then $\left\{r_{i} \pm s_{i}\right\} \sim\left\{r_{i}^{\prime} \pm s_{i}^{\prime}\right\}$ and $\left\{r_{i} \cdot s_{i}\right\} \sim\left\{r_{i}^{\prime} \cdot s_{i}^{\prime}\right\}$. Similarly, $\left\{r_{i} / s_{i}\right\} \sim\left\{r_{i}^{\prime} / s_{i}^{\prime}\right\}$ provided $\left\{s_{i}\right\} \nsim \overline{0}$, and $s_{i} \neq 0$ and $s_{i}^{\prime} \neq 0$ for every $i$.
2.13. Definition. If $\rho$ and $\sigma$ are represented by $\left\{r_{i}\right\}$ and $\left\{s_{i}\right\}$ respectively, then $\rho \pm \sigma$ is defined by the equivalence class containing $\left\{r_{i} \pm s_{i}\right\}$ and $\rho \cdot \sigma$ by $\left\{r_{i} \cdot s_{i}\right\}$. $\rho / \sigma$ is defined to be the equivalence class containing $\left\{r_{i} / s_{i}^{\prime}\right\}$ where $\left\{s_{i}\right\} \sim\left\{s_{i}^{\prime}\right\}$ and $s_{i}^{\prime} \neq 0$ for all $i$, provided $\left\{s_{i}\right\} \nsim \overline{0}$.

Reference to Propositions 2.11 and 2.12 shows that these operations are well defined. That is, if $\rho^{\prime}$ and $\sigma^{\prime}$ are represented by $\left\{r_{i}^{\prime}\right\}$ and $\left\{s_{i}^{\prime}\right\}$, where $\left\{r_{i}\right\} \sim\left\{r_{i}^{\prime}\right\}$ and $\left\{s_{i}\right\} \sim\left\{s_{i}^{\prime}\right\}$, then $\rho+\sigma=\rho^{\prime}+\sigma^{\prime}$ and similarly for the other operations.

Since the rational numbers form a field, it is clear that the real numbers also form a field. However, we wish to show that they actually form an Archimedean ordered field. For this we first must define an ordering on the real numbers that is compatible with the field structure; this will be accomplished by the following theorem.
2.14. Theorem. If $\left\{r_{i}\right\}$ and $\left\{s_{i}\right\}$ are Cauchy, then one (and only one) of the following occurs:
(i) $\left\{r_{i}\right\} \sim\left\{s_{i}\right\}$.
(ii) There exist a positive integer $N$ and a positive rational number $k$ such that $r_{i}>s_{i}+k$ for $i \geq N$.
(iii) There exist a positive integer $N$ and positive rational number $k$ such that $s_{i}>r_{i}+k$ for $i \geq N$.

Proof. Suppose that (i) does not hold. Then there exists a rational number $k>0$ with the property that for every positive integer $N$ there exists an integer $i \geq N$ such that

$$
\left|r_{i}-s_{i}\right|>2 k
$$

This is equivalent to saying that

$$
\left|r_{i}-s_{i}\right|>2 k \quad \text { for infinitely many } \quad i \geq 1
$$

Since $\left\{r_{i}\right\}$ is Cauchy, there exists a positive integer $N$ such that

$$
\left|r_{i}-r_{j}\right|<k / 2 \quad \text { for all } \quad i, j \geq N_{1}
$$

Likewise, there exists a positive integer $N_{2}$ such that

$$
\left|s_{i}-s_{j}\right|<k / 2 \quad \text { for all } \quad i, j \geq N_{2}
$$

Let $N^{*} \geq \max \left\{N_{1}, N_{2}\right\}$ be an integer with the property that

$$
\left|r_{N^{*}}-s_{N^{*}}\right|>2 k
$$

Either $r_{N^{*}}>s_{N^{*}}$ or $s_{N^{*}}>r_{N^{*}}$. We will show that the first possibility leads to conclusion (ii) of the theorem. The proof that the second possibility leads to (iii) is similar and will be omitted. Assuming now that $r_{N^{*}}>s_{N^{*}}$, we have

$$
r_{N^{*}}>s_{N^{*}}+2 k
$$

It follows from (2.1) and (2.4) that

$$
\left|r_{N^{*}}-r_{i}\right|<k / 2 \quad \text { and } \quad\left|s_{N^{*}}-s_{i}\right|<k / 2 \quad \text { for all } \quad i \geq N^{*}
$$

From this and (2.6) we have that

$$
r_{i}>r_{N^{*}}-k / 2>s_{N^{*}}+2 k-k / 2=s_{N^{*}}+3 k / 2 \quad \text { for } \quad i \geq N^{*}
$$

But $s_{N^{*}}>s_{i}-k / 2$ for $i \geq N^{*}$ and consequently,

$$
r_{i}>s_{i}+k \quad \text { for } \quad i \geq N^{*} .
$$

2.15. Definition. If $\rho=\overline{\left\{r_{i}\right\}}$ and $\sigma=\overline{\left\{s_{i}\right\}}$, then we say that $\rho<\sigma$ if there exist rational numbers $q_{1}$ and $q_{2}$ with $q_{1}<q_{2}$ and a positive integer $N$ such that such that $r_{i}<q_{1}<q_{2}<s_{i}$ for all $i$ with $i \geq N$. Note that $q_{1}$ and $q_{2}$ can be chosen to be independent of the representative Cauchy sequences of rational numbers that determine $\rho$ and $\sigma$.

In view of this definition, Theorem 2.14 implies that the real numbers are comparable, which we state in the following corollary.
2.16. Theorem. Corollary If $\rho$ and $\sigma$ are real numbers, then one (and only one) of the following must hold:
(1) $\rho=\sigma$,
(2) $\rho<\sigma$,
(3) $\rho>\sigma$.

Moreover, $\mathbb{R}$ is an Archimedean ordered field.

The compatibility of $\leq$ with the field structure of $\mathbb{R}$ follows from Theorem 2.14. That $\mathbb{R}$ is Archimedean follows from Theorem 2.14 and the fact that $\mathbb{Q}$ is Archimedean. Note that the absolute value of a real number can thus be defined analogously to that of a rational number.
2.17. Definition. If $\left\{\rho_{i}\right\}_{i=1}^{\infty}$ is a sequence in $\mathbb{R}$ and $\rho \in \mathbb{R}$ we define

$$
\lim _{i \rightarrow \infty} \rho_{i}=\rho
$$

to mean that given any real number $\varepsilon>0$ there is a positive integer $N$ such that

$$
\left|\rho_{i}-\rho\right|<\varepsilon \quad \text { whenever } \quad i \geq N
$$

2.18. Remark. Having shown that $\mathbb{R}$ is an Archimedean ordered field, we now know that $\mathbb{Q}$ has a natural injection into $\mathbb{R}$ by way of the constant sequences. That is, if $r \in \mathbb{Q}$, then the constant sequence $\bar{r}$ gives its corresponding equivalence class in $\mathbb{R}$. Consequently, we shall consider $\mathbb{Q}$ to be a subset of $\mathbb{R}$, that is, we do not distinguish between $r$ and its corresponding equivalence class. Moreover, if $\rho_{1}$ and $\rho_{2}$ are in $\mathbb{R}$ with $\rho_{1}<\rho_{2}$, then there is a rational number $r$ such that $\rho_{1}<r<\rho_{2}$.

The next proposition provides a connection between Cauchy sequences in $\mathbb{Q}$ with convergent sequences in $\mathbb{R}$.
2.19. Theorem. If $\rho=\overline{\left\{r_{i}\right\}}$, then

$$
\lim _{i \rightarrow \infty} r_{i}=\rho
$$

Proof. Given $\varepsilon>0$, we must show the existence of a positive integer $N$ such that $\left|r_{i}-\rho\right|<\varepsilon$ whenever $i \geq N$. Let $\varepsilon$ be represented by the rational sequence $\left\{\varepsilon_{i}\right\}$. Since $\varepsilon>0$, we know from Theorem (2.14), (ii), that there exist a positive rational number $k$ and an integer $N_{1}$ such that $\varepsilon_{i}>k$ for all $i \geq N_{1}$. Because the sequence $\left\{r_{i}\right\}$ is Cauchy, we know there exists a positive integer $N_{2}$ such that $\left|r_{i}-r_{j}\right|<k / 2$ whenever $i, j \geq N_{2}$. Fix an integer $i \geq N_{2}$ and let $r_{i}$ be determined by the constant sequence $\left\{r_{i}, r_{i}, \ldots\right\}$. Then the real number $\rho-r_{i}$ is determined by the Cauchy sequence $\left\{r_{j}-r_{i}\right\}$, that is

$$
\rho-r_{i}=\overline{\left\{r_{j}-r_{i}\right\}}
$$

If $j \geq N_{2}$, then $\left|r_{j}-r_{i}\right|<k / 2$. Note that the real number $\left|\rho-r_{i}\right|$ is determined by the sequence $\left\{\left|r_{j}-r_{i}\right|\right\}$. Now, the sequence $\left\{\left|r_{j}-r_{i}\right|\right\}$ has the property that $\left|r_{j}-r_{i}\right|<k / 2<k<\varepsilon_{j}$ for $j \geq \max \left(N_{1}, N_{2}\right)$. Hence, by Definition (2.15), $\left|\rho-r_{i}\right|<\varepsilon$. The proof is concluded by taking $N=\max \left(N_{1}, N_{2}\right)$.
2.20. Theorem. The set of real numbers is complete; that is, every Cauchy sequence of real numbers converges to a real number.

Proof. Let $\left\{\rho_{i}\right\}$ be a Cauchy sequence of real numbers and let each $\rho_{i}$ be determined by the Cauchy sequence of rational numbers, $\left\{r_{i, k}\right\}_{k=1}^{\infty}$. By the previous proposition,

$$
\lim _{k \rightarrow \infty} r_{i, k}=\rho_{i}
$$

Thus, for each positive integer $i$, there exists $k_{i}$ such that

$$
\begin{equation*}
\left|r_{i, k_{i}}-\rho_{i}\right|<\frac{1}{i} \tag{2.1}
\end{equation*}
$$

Let $s_{i}=r_{i, k_{i}}$. The sequence $\left\{s_{i}\right\}$ is Cauchy because

$$
\begin{aligned}
\left|s_{i}-s_{j}\right| & \leq\left|s_{i}-\rho_{i}\right|+\left|\rho_{i}-\rho_{j}\right|+\left|\rho_{j}-s_{j}\right| \\
& \leq 1 / i+\left|\rho_{i}-\rho_{j}\right|+1 / j .
\end{aligned}
$$

Indeed, for $\varepsilon>0$, there exists a positive integer $N>4 / \varepsilon$ such that $i, j \geq N$ implies $\left|\rho_{i}-\rho_{j}\right|<\varepsilon / 2$. This, along with (2.1), shows that $\left|s_{i}-s_{j}\right|<\varepsilon$ for $i, j \geq N$. Moreover, if $\rho$ is the real number determined by $\left\{s_{i}\right\}$, then

$$
\begin{aligned}
\left|\rho-\rho_{i}\right| & \leq\left|\rho-s_{i}\right|+\left|s_{i}-\rho_{i}\right| \\
& \leq\left|\rho-s_{i}\right|+1 / i
\end{aligned}
$$

For $\varepsilon>0$, we invoke Theorem 2.19 for the existence of $N>2 / \varepsilon$ such that the first term is less than $\varepsilon / 2$ for $i \geq N$. For all such $i$, the second term is also less than $\varepsilon / 2$.

The completeness of the real numbers leads to another property that is of basic importance.
2.21. Definition. A number $M$ is called an upper bound for for a set $A \subset \mathbb{R}$ if $a \leq M$ for all $a \in A$. An upper bound $b$ for $A$ is called a least upper bound for $A$ if $b$ is less than all other upper bounds for $A$. The term supremum of $A$ is used interchangeably with least upper bound and is written sup $A$. The terms lower bound, greatest lower bound, and infimum are defined analogously.
2.22. Theorem. Let $A \subset \mathbb{R}$ be a nonempty set that is bounded above (below). Then $\sup A(\inf A)$ exists.

Proof. Let $b \in \mathbb{R}$ be any upper bound for $A$ and let $a \in A$ be an arbitrary element. Further, using the Archimedean property of $\mathbb{R}$, let $M$ and $-m$ be positive integers such that $M>b$ and $-m>-a$, so that we have $m<a \leq b<M$. For each positive integer $p$ let

$$
I_{p}=\left\{k: k \text { an integer and } \frac{k}{2^{p}} \text { is an upper bound for } A\right\} .
$$

Since $A$ is bounded above, it follows that $I_{p}$ is not empty. Furthermore, if $a \in A$ is an arbitrary element, there is an integer $j$ that is less than $a$. If $k$ is an integer such that $k \leq 2^{p} j$, then $k$ is not an element of $I_{p}$, thus showing that $I_{p}$ is bounded
below. Therefore, since $I_{p}$ consists only of integers, it follows that $I_{p}$ has a first element, call it $k_{p}$. Because

$$
\frac{2 k_{p}}{2^{p+1}}=\frac{k_{p}}{2^{p}}
$$

the definition of $k_{p+1}$ implies that $k_{p+1} \leq 2 k_{p}$. But

$$
\frac{2 k_{p}-2}{2^{p+1}}=\frac{k_{p}-1}{2^{p}}
$$

is not an upper bound for $A$, which implies that $k_{p+1} \neq 2 k_{p}-2$. In fact, it follows that $k_{p+1}>2 k_{p}-2$. Therefore, either

$$
k_{p+1}=2 k_{p} \quad \text { or } \quad k_{p+1}=2 k_{p}-1
$$

Defining $a_{p}=\frac{k_{p}}{2^{p}}$, we have either

$$
a_{p+1}=\frac{2 k_{p}}{2^{p+1}}=a_{p} \quad \text { or } \quad a_{p+1}=\frac{2 k_{p}-1}{2^{p+1}}=a_{p}-\frac{1}{2^{p+1}},
$$

and hence,

$$
a_{p+1} \leq a_{p} \quad \text { with } \quad a_{p}-a_{p+1} \leq \frac{1}{2^{p+1}}
$$

for each positive integer $p$. If $q>p \geq 1$, then

$$
\begin{aligned}
0 \leq a_{p}-a_{q} & =\left(a_{p}-a_{p+1}\right)+\left(a_{p+1}-a_{p+2}\right)+\cdots+\left(a_{q-1}-a_{q}\right) \\
& \leq \frac{1}{2^{p+1}}+\frac{1}{2^{p+2}}+\cdots+\frac{1}{2^{q}} \\
& =\frac{1}{2^{p+1}}\left(1+\frac{1}{2}+\cdots+\frac{1}{2^{q-p-1}}\right) \\
& <\frac{1}{2^{p+1}}(2)=\frac{1}{2^{p}}
\end{aligned}
$$

Thus, whenever $q>p \geq 1$, we have $\left|a_{p}-a_{q}\right|<\frac{1}{2^{p}}$, which implies that $\left\{a_{p}\right\}$ is a Cauchy sequence. By the completeness of the real numbers, Theorem 2.20, there is a real number $c$ to which the sequence converges.

We will show that $c$ is the supremum of $A$. First, observe that $c$ is an upper bound for $A$ since it is the limit of a decreasing sequence of upper bounds. Secondly, it must be the least upper bound, for otherwise, there would be an upper bound $c^{\prime}$ with $c^{\prime}<c$. Choose an integer $p$ such that $1 / 2^{p}<c-c^{\prime}$. Then

$$
a_{p}-\frac{1}{2^{p}} \geq c-\frac{1}{2^{p}}>c+c^{\prime}-c=c^{\prime}
$$

which shows that $a_{p}-\frac{1}{2^{p}}$ is an upper bound for $A$. But the definition of $a_{p}$ implies that

$$
a_{p}-\frac{1}{2^{p}}=\frac{k_{p}-1}{2^{p}}
$$

a contradiction, since $\frac{k_{p}-1}{2^{p}}$ is not an upper bound for $A$.
The existence of $\inf A$ in case $A$ is bounded below follows by analogous argument.

A linearly ordered field is said to have the least upper bound property if each nonempty subset that has an upper bound has a least upper bound (in the field). Hence, $\mathbb{R}$ has the least upper bound property. It can be shown that every linearly ordered field with the least upper bound property is a complete Archimedean ordered field. We will not prove this assertion.

## Exercises for Section 2.1

1. Use the fact that

$$
\mathbb{N}=\{n: n=2 k \text { for some } k \in \mathbb{N}\} \cup\{n: n=2 k+1 \text { for some } k \in \mathbb{N}\}
$$

to prove $c \cdot c=c$. Consequently, $\operatorname{card}\left(\mathbb{R}^{n}\right)=c$ for each $n \in \mathbb{N}$.
2. Suppose $\alpha, \beta$ and $\delta$ are cardinal numbers. Prove that

$$
\delta^{\alpha+\beta}=\delta^{\alpha} \cdot \delta^{\beta}
$$

3. Prove that the set of numbers whose dyadic expansions are not unique is countable.
4. Prove that the equation $x^{2}-2=0$ has no solutions in the field $\mathbb{Q}$.
5. Prove: If $\left\{x_{n}\right\}_{n=1}^{\infty}$ is a bounded, increasing sequence in an Archimedean ordered field, then the sequence is Cauchy.
6. Prove that each Archimedean ordered field contains a "copy" of $\mathbb{Q}$. Moreover, for each pair $r_{1}$ and $r_{2}$ of the field with $r_{1}<r_{2}$, there exists a rational number $r$ such that $r_{1}<r<r_{2}$.
7. Consider the set $\{r+q \sqrt{2}: r \in \mathbb{Q}, q \in \mathbb{Q}\}$. Prove that it is an Archimedean ordered field.
8. Let $F$ be the field of all rational polynomials with coefficients in $\mathbb{Q}$. Thus, a typical element of $F$ has the form $\frac{P(x)}{Q(x)}$, where $P(x)=\sum_{k=0}^{n} a_{k} x^{k}$ and $Q(x)=$ $\sum_{j=0}^{m} b_{j} x^{j}$ where the $a_{k}$ and $b_{j}$ are in $\mathbb{Q}$ with $a_{n} \neq 0$ and $b_{m} \neq 0$. We order $F$ by saying that $\frac{P(x)}{Q(x)}$ is positive if and only if $a_{n} b_{m}$ is a positive rational number. Prove that $F$ is an ordered field which is not Archimedean.
9. Consider the set $\{0,1\}$ with + and $\times$ given by the following tables:

| + | 0 | 1 |
| :---: | :---: | :---: |
| 0 | 0 | 1 |
| 1 | 1 | 0 |


| $\times$ | 0 | 1 |
| :---: | :---: | :---: |
| 0 | 0 | 0 |
| 1 | 0 | 1 |

Prove that $\{0,1\}$ is a field and that there can be no ordering on $\{0,1\}$ that results in a linearly ordered field.
10. Prove: For real numbers $a$ and $b$,
(a) $|a+b| \leq|a|+|b|$,
(b) $||a|-|b|| \leq|a-b|$
(c) $|a b|=|a||b|$.

### 2.2. Cardinal Numbers

There are many ways to determine the "size" of a set, the most basic being the enumeration of its elements when the set is finite. When the set is infinite, another means must be employed; the one that we use is not far from the enumeration concept.
2.23. Definition. Two sets $A$ and $B$ are said to be equivalent if there exists a bijection $f: A \rightarrow B$, and then we write $A \sim B$. In other words, there is a one-to-one correspondence between $A$ and $B$. It is not difficult to show that this notion of equivalence defines an equivalence relation as described in Definition 1.1 and therefore sets are partitioned into equivalence classes. Two sets in the same equivalence class are said to have the same cardinal number or to be of the same cardinality. The cardinal number of a set $A$ is denoted by card $A$; that is, card $A$ is the symbol we attach to the equivalence class containing $A$. There are some sets so frequently encountered that we use special symbols for their cardinal numbers. For example, the cardinal number of the set $\{1,2, \ldots, n\}$ is denoted by $n, \operatorname{card} \mathbb{N}=\aleph_{0}$, and $\operatorname{card} \mathbb{R}=c$.
2.24. Definition. Let $A$ be a non-empty set. If card $A=n$, for some nonnegative integer $n$, then we say that $A$ is a finite set. If $A$ is not finite then we say that it is an infinite set. If $A$ is equivalent to the positive integers then $A$ is denumerable. If $A$ is either finite or denumerable then it is called countable; otherwise it is called uncountable.

One of the first observations concerning cardinality is that it is possible for two sets to have the same cardinality even though one is a proper subset of the other. For example, the formula $y=2 x, x \in[0,1]$ defines a bijection between the closed intervals $[0,1]$ and $[0,2]$. This also can be seen with the help of the figure below.


Another example, utilizing a two-step process, establishes the equivalence between points $x$ of $(-1,1)$ and $y$ of $\mathbb{R}$. The semicircle with endpoints omitted serves as an intermediary.


A bijection could also be explicitly given by $y=\frac{2 x-1}{1-(2 x-1)^{2}}, x \in(0,1)$.
Pursuing other examples, it should be true that $(0,1) \sim[0,1]$ although in this case, exhibiting the bijection is not immediately obvious (but not very difficult, see Exercise 7 at the end of this section). Aside from actually exhibiting the bijection, the facts that $(0,1)$ is equivalent to a subset of $[0,1]$ and that $[0,1]$ is equivalent to a subset of $(0,1)$ offer compelling evidence that $(0,1) \sim[0,1]$. The next two results make this rigorous.
2.25. Theorem. If $A \supset A_{1} \supset A_{2}$ and $A \sim A_{2}$, then $A \sim A_{1}$.

Proof. Let $f: A \rightarrow A_{2}$ denote the bijection that determines the equivalence between $A$ and $A_{2}$. The restriction of $f$ to $A_{1}, f\left\llcorner A_{1}\right.$, determines a set $A_{3}$ (actually, $\left.A_{3}=f\left(A_{1}\right)\right)$ such that $A_{1} \sim A_{3}$ where $A_{3} \subset A_{2}$. Now we have sets $A_{1} \supset A_{2} \supset A_{3}$ such that $A_{1} \sim A_{3}$. Repeating the argument, there exists a set $A_{4}, A_{4} \subset A_{3}$ such that $A_{2} \sim A_{4}$. Continue this way to obtain a sequence of sets such that

$$
A \sim A_{2} \sim A_{4} \sim \cdots \sim A_{2 i} \sim \cdots
$$

and

$$
A_{1} \sim A_{3} \sim A_{5} \sim \cdots \sim A_{2 i+1} \sim \cdots
$$

For notational convenience, we take $A_{0}=A$. Then we have

$$
\begin{gather*}
A_{0}=\left(A_{0}-A_{1}\right) \cup\left(A_{1}-A_{2}\right) \cup\left(A_{2}-A_{3}\right) \cup \cdots  \tag{2.2}\\
\cup\left(A_{0} \cap A_{1} \cap A_{2} \cap \cdots\right)
\end{gather*}
$$

and

$$
\begin{align*}
& A_{1}=\left(A_{1}-A_{2}\right) \cup\left(A_{2}-A_{3}\right) \cup\left(A_{3}-A_{4}\right) \cup \cdots  \tag{2.3}\\
& \cup\left(A_{1} \cap A_{2} \cap A_{3} \cap \cdots\right)
\end{align*}
$$

By the properties of the sets constructed, we see that

$$
\begin{equation*}
\left(A_{0}-A_{1}\right) \sim\left(A_{2}-A_{3}\right),\left(A_{2}-A_{3}\right) \sim\left(A_{4}-A_{5}\right), \cdots \tag{2.4}
\end{equation*}
$$

In fact, the bijection between $\left(A_{0}-A_{1}\right)$ and $\left(A_{2}-A_{3}\right)$ is given by $f$ restricted to $A_{0}-A_{1}$. Likewise, $f$ restricted to $A_{2}-A_{3}$ provides a bijection onto $A_{4}-A_{5}$, and similarly for the remaining sets in the sequence. Moreover, since $A_{0} \supset A_{1} \supset A_{2} \supset$ $\cdots$, we have

$$
\left(A_{0} \cap A_{1} \cap A_{2} \cap \cdots\right)=\left(A_{1} \cap A_{2} \cap A_{3} \cap \cdots\right) .
$$

The sets $A_{0}$ and $A_{1}$ are represented by a disjoint union of sets in (2.2) and (2.3). With the help of (2.4), note that the union of the first two sets that appear in the expressions for $A$ and in $A_{1}$ are equivalent; that is,

$$
\left(A_{0}-A_{1}\right) \cup\left(A_{1}-A_{2}\right) \sim\left(A_{1}-A_{2}\right) \cup\left(A_{2}-A_{3}\right) .
$$

Likewise,

$$
\left(A_{2}-A_{3}\right) \cup\left(A_{4}-A_{5}\right) \sim\left(A_{3}-A_{4}\right) \cup\left(A_{5}-A_{6}\right),
$$

and similarly for the remaining sets. Thus, it is easy to see that $A \sim A_{1}$.
2.26. Theorem (Schröder-Bernstein). If $A \supset A_{1}, B \supset B_{1}, A \sim B_{1}$ and $B \sim A_{1}$, then $A \sim B$.

Proof. Denoting by $f$ the bijection that determines the similarity between $A$ and $B_{1}$, let $B_{2}=f\left(A_{1}\right)$ to obtain $A_{1} \sim B_{2}$ with $B_{2} \subset B_{1}$. However, by hypothesis, we have $A_{1} \sim B$ and therefore $B \sim B_{2}$. Now invoke Theorem 2.25 to conclude that $B \sim B_{1}$. But $A \sim B_{1}$ by hypothesis and consequently, $A \sim B$.

It is instructive to recast all of the information in this section in terms of cardinality. First, we introduce the concept of comparability of cardinal numbers.
2.27. Definition. If $\alpha$ and $\beta$ are the cardinal numbers of the sets A and B , respectively, we say $\alpha \leq \beta$ if and only if there exists a set $B_{1} \subset B$ such that $A \sim B_{1}$. In addition, we say that $\alpha<\beta$ if there exists no set $A_{1} \subset A$ such that $A_{1} \sim B$.

With this terminology, the Schröder-Bernstein Theorem states that

$$
\alpha \leq \beta \quad \text { and } \quad \beta \leq \alpha \quad \text { implies } \quad \alpha=\beta .
$$

The next definition introduces arithmetic operations on the cardinal numbers.
2.28. Definition. Using the notation of Definition 2.27 we define

$$
\begin{aligned}
\alpha+\beta & =\operatorname{card}(A \cup B) \quad \text { where } \quad A \cap B=\emptyset \\
\alpha \cdot \beta & =\operatorname{card}(A \times B) \\
\alpha^{\beta} & =\operatorname{card} F
\end{aligned}
$$

where $F$ is the family of all functions $f: B \rightarrow A$.
Let us examine the last definition in the special case where $\alpha=2$. If we take the corresponding set $A$ as $A=\{0,1\}$, it is easy to see that $F$ is equivalent to the class of all subsets of $B$. Indeed, the bijection can be defined by

$$
f \rightarrow f^{-1}\{1\}
$$

where $f \in F$. This bijection is nothing more than correspondence between subsets of $B$ and their associated characteristic functions. Thus, $2^{\beta}$ is the cardinality of all subsets of $B$, which agrees with what we already know in case $\beta$ is finite. Also, from previous discussions in this section, we have

$$
\aleph_{0}+\aleph_{0}=\aleph_{0}, \aleph_{0} \cdot \aleph_{0}=\aleph_{0} \quad \text { and } \quad c+c=c
$$

In addition, we see that the customary basic arithmetic properties are preserved.
2.29. Theorem. If $\alpha, \beta$ and $\gamma$ are cardinal numbers, then
(i) $\alpha+(\beta+\gamma)=(\alpha+\beta)+\gamma$
(ii) $\alpha(\beta \gamma)=(\alpha \beta) \gamma$
(iii) $\alpha+\beta=\beta+\alpha$
(iv) $\alpha^{(\beta+\gamma)}=\alpha^{\beta} \alpha^{\gamma}$
(v) $\alpha^{\gamma} \beta^{\gamma}=(\alpha \beta)^{\gamma}$
(vi) $\left(\alpha^{\beta}\right)^{\gamma}=\alpha^{\beta \gamma}$

The proofs of these properties are quite easy. We give an example by proving (vi):

Proof of (Vi). Assume that sets $A, B$ and $C$ respectively represent the cardinal numbers $\alpha, \beta$ and $\gamma$. Recall that $\left(\alpha^{\beta}\right)^{\gamma}$ is represented by the family $\mathcal{F}$ of all mappings $f$ defined on $C$ where $f(c): B \rightarrow A$. Thus, $f(c)(b) \in A$. On the other hand, $\alpha^{\beta \gamma}$ is represented by the family $\mathcal{G}$ of all mappings $g: B \times C \rightarrow A$. Define

$$
\varphi: \mathcal{F} \rightarrow \mathcal{G}
$$

as $\varphi(f)=g$ where

$$
g(b, c):=f(c)(b)
$$

that is,

$$
\varphi(f)(b, c)=f(c)(b)=g(b, c)
$$

Clearly, $\varphi$ is surjective. To show that $\varphi$ is univalent, let $f_{1}, f_{2} \in \mathcal{F}$ be such that $f_{1} \neq f_{2}$. For this to be true, there exists $c_{0} \in C$ such that

$$
f_{1}\left(c_{0}\right) \neq f_{2}\left(c_{0}\right)
$$

This, in turn, implies the existence of $b_{0} \in B$ such that

$$
f_{1}\left(c_{0}\right)\left(b_{0}\right) \neq f_{2}\left(c_{0}\right)\left(b_{0}\right)
$$

and this means that $\varphi\left(f_{1}\right)$ and $\varphi\left(f_{2}\right)$ are different mappings, as desired.
In addition to these arithmetic identities, we have the following theorems which deserve special attention.
2.30. Theorem. $2^{\aleph_{0}}=c$.

Proof. First, to prove the inequality $2^{\aleph_{0}} \geq c$, observe that each real number $r$ is uniquely associated with the subset $Q_{r}:=\{q: q \in \mathbb{Q}, q<r\}$ of $\mathbb{Q}$. Thus mapping $r \mapsto Q_{r}$ is an injection from $\mathbb{R}$ into $\mathcal{P}(\mathbb{Q})$. Hence,

$$
c=\operatorname{card} \mathbb{R} \leq \operatorname{card}[\mathcal{P}(\mathbb{Q})]=\operatorname{card}[\mathcal{P}(\mathbb{N})]=2^{\aleph_{0}}
$$

because $\mathbb{Q} \sim \mathbb{N}$.
To prove the opposite inequality, consider the set $\mathcal{S}$ of all sequences of the form $\left\{x_{k}\right\}$ where $x_{k}$ is either 0 or 1 . Referring to the definition of a sequence (Definition 1.2 ), it is immediate that the cardinality of $\mathcal{S}$ is $2^{\aleph_{0}}$. We will see below (Corollary 2.36) that each number $x \in[0,1]$ has a decimal representation of the form

$$
x=. x_{1} x_{2} \ldots, x_{i} \in\{0,1\}
$$

Of course, such representations do not uniquely represent $x$. For example,

$$
\frac{1}{2}=.10000 \ldots=.01111 \ldots
$$

Accordingly, the mapping from $\mathcal{S}$ into $\mathbb{R}$ defined by

$$
f\left(\left\{x_{k}\right\}\right)= \begin{cases}\sum_{k=1}^{\infty} \frac{x_{k}}{2^{k}} & \text { if } x_{k} \neq 0 \text { for all but finitely many } k \\ \sum_{k=1}^{\infty} \frac{x_{k}}{2^{k}}+1 & \text { if } x_{k}=0 \text { for infinitely many } k\end{cases}
$$

is clearly an injection, thus proving that $2^{\aleph_{0}} \leq c$. Now apply the Schröder-Bernstein Theorem to obtain our result.

The previous result implies, in particular, that $2^{\aleph_{0}}>\aleph_{0}$; the next result is a generalization of this.
2.31. Theorem. For any cardinal number $\alpha, 2^{\alpha}>\alpha$.

Proof. If $A$ has cardinal number $\alpha$, it follows that $2^{\alpha} \geq \alpha$ since each element of $A$ determines a singleton that is a subset of $A$. Proceeding by contradiction, suppose $2^{\alpha}=\alpha$. Then there exists a one-to-one correspondence between elements $x$ and sets $S_{x}$ where $x \in A$ and $S_{x} \subset A$. Let $D=\left\{x \in A: x \notin S_{x}\right\}$. By assumption there exists $x_{0} \in A$ such that $x_{0}$ is related to the set $D$ under the one-to-one correspondence (i.e., $D=S_{x_{0}}$ ). However, this leads to a contradiction; consider the following two possibilities:
(1) If $x_{0} \in D$, then $x_{0} \notin S_{x_{0}}$ by the definition of $D$. But then, $x_{0} \notin D$, a contradiction.
(2) If $x_{0} \notin D$, similar reasoning leads to the conclusion that $x_{0} \in D$

The next proposition, whose proof is left to the reader, shows that $\aleph_{0}$ is the smallest infinite cardinal.
2.32. Proposition. Every infinite set $S$ contains a denumerable subset.

An immediate consequence of the proposition is the following characterization of infinite sets.
2.33. Theorem. A nonempty set $S$ is infinite if and only if for each $x \in S$ the sets $S$ and $S-\{x\}$ are equivalent

By means of the Schröder-Berstein theorem, it is now easy to show that the rationals are denumerable. In fact, we show a bit more.
2.34. Proposition. (i) The set of rational numbers is denumerable, (ii) If $A_{i}$ is denumerable for $i \in \mathbb{N}$, then $A:=\bigcup_{i \in \mathbb{N}} A_{i}$ is denumerable.

Proof. Case (i) is subsumed by (ii). Since the sets $A_{i}$ are denumerable, their elements can be enumerated by $\left\{a_{i, 1}, a_{i, 2}, \ldots\right\}$. For each $a \in A$, let $\left(k_{a}, j_{a}\right)$ be the unique pair in $\mathbb{N} \times \mathbb{N}$ such that

$$
k_{a}=\min \left\{k: a=a_{k, j}\right\}
$$

and

$$
j_{a}=\min \left\{j: a=a_{k_{a}, j}\right\}
$$

(Be aware that $a$ could be present more than once in $A$. If we visualize $A$ as an infinite matrix, then $\left(k_{a}, j_{a}\right)$ represents the position of $a$ that is furthest to the
"northwest" in the matrix.) Consequently, $A$ is equivalent to a subset of $\mathbb{N} \times \mathbb{N}$. Further, observe that there is an injection of $\mathbb{N} \times \mathbb{N}$ into $\mathbb{N}$ given by

$$
(i, j) \rightarrow 2^{i} 3^{j}
$$

Indeed, if this were not an injection, we would have

$$
2^{i-i^{\prime}} 3^{j-j^{\prime}}=1
$$

for some distinct positive integers $i, i^{\prime}, j$, and $j^{\prime}$, which is impossible. Thus, it follows that $A$ is equivalent to a subset of $\mathbb{N}$ and is therefore equivalent to a subset of $A_{1}$ because $\mathbb{N} \sim A_{1}$. Since $A_{1} \subset A$ we can now appeal to the Schröder-Bernstein Theorem to arrive at the desired conclusion.

It is natural to ask whether the real numbers are also denumerable. This turns out to be false, as the following two results indicate. It was G. Cantor who first proved this fact.
2.35. Theorem. If $I_{1} \supset I_{2} \supset I_{3} \supset \ldots$ are closed intervals with the property that length $I_{i} \rightarrow 0$, then

$$
\bigcap_{i=1}^{\infty} I_{i}=\left\{x_{0}\right\}
$$

for some point $x_{0} \in \mathbb{R}$.
Proof. Let $I_{i}=\left[a_{i}, b_{i}\right]$ and choose $x_{i} \in I_{i}$. Then $\left\{x_{i}\right\}$ is a Cauchy sequence of real numbers since $\left|x_{i}-x_{j}\right| \leq \max \left[\operatorname{length} I_{i}\right.$, length $\left.I_{j}\right]$. Since $\mathbb{R}$ is complete (Theorem 2.20), there exists $x_{0} \in \mathbb{R}$ such that

$$
\begin{equation*}
\lim _{i \rightarrow \infty} x_{i}=x_{0} \tag{2.5}
\end{equation*}
$$

We claim that

$$
\begin{equation*}
x_{0} \in \bigcap_{i=1}^{\infty} I_{i} \tag{2.6}
\end{equation*}
$$

for if not, there would be some positive integer $i_{0}$ for which $x_{0} \notin I_{i_{0}}$. Therefore, since $I_{i_{0}}$ is closed, there would be an $\eta>0$ such that $\left|x_{0}-y\right|>\eta$ for each $y \in I_{i_{0}}$. Since the intervals are nested, it would follow that $x_{0} \notin I_{i}$ for all $i \geq i_{0}$ and thus $\left|x_{0}-x_{i}\right|>\eta$ for all $i \geq i_{0}$. This would contradict (2.5) thus establishing (2.6). We leave it to the reader to verify that $x_{0}$ is the only point with this property.
2.36. Corollary. Every real number has a decimal representation relative to any base.
2.37. Theorem. The real numbers are uncountable.

Proof. The proof proceeds by contradiction. Thus, we assume that the real numbers can be enumerated as $a_{1}, a_{2}, \ldots, a_{i}, \ldots$ Let $I_{1}$ be a closed interval of positive length less than 1 such that $a_{1} \notin I_{1}$. Let $I_{2} \subset I_{1}$ be a closed interval of positive length less than $1 / 2$ such that $a_{2} \notin I_{2}$. Continue in this way to produce a nested sequence of intervals $\left\{I_{i}\right\}$ of positive length than $1 / i$ with $a_{i} \notin I_{i}$. Lemma 2.35 , we have the existence of a point

$$
x_{0} \in \bigcap_{i=1}^{\infty} I_{i} .
$$

Observe that $x_{0} \neq a_{i}$ for any $i$, contradicting the assumption that all real numbers are among the $a_{i}$ 's.

## Exercises for Section 2.2

1. Show that an arbitrary function $\mathbb{R} \xrightarrow{f} \mathbb{R}$ has at most a countable number of removable discontinuities: that is, prove that

$$
A:=\left\{a \in \mathbb{R}: \lim _{x \rightarrow a} f(x) \text { exists and } \lim _{x \rightarrow a} f(x) \neq f(a)\right\}
$$

is at most countable.
2. Show that an arbitrary function $\mathbb{R} \xrightarrow{f} \mathbb{R}$ has at most a countable number of jump discontinuities: that is, let

$$
f^{+}(a):=\lim _{x \rightarrow a^{+}} f(x)
$$

and

$$
f^{-}(a):=\lim _{x \rightarrow a^{-}} f(x) .
$$

Show that the set $\left\{a \in \mathbb{R}: f^{+}(a) \neq f^{-}(a)\right\}$ is at most countable.
3. Prove: If $A$ is the union of a countable collection of countable sets, then $A$ is a countable set.
4. Prove Proposition 2.33.
5. Let $B$ be a countable subset of an uncountable set $A$. Show that $A$ is equivalent to $A \backslash B$.
6. Prove that a set $A \subset \mathbb{N}$ is finite if and only if $A$ has an upper bound.
7. Exhibit an explicit bijection between $(0,1)$ and $[0,1]$.
8. If you are working in Zermelo-Fraenkel set theory without the Axiom of Choice, can you choose an element from...
a finite set?
an infinite set?
each member of an infinite set of singletons (i.e., one-element sets)?
each member of an infinite set of pairs of shoes?
each member of an infinite set of pairs of socks?
each member of a finite set of sets if each of the members is infinite?
each member of an infinite set of sets if each of the members is infinite?
each member of a denumerable set of sets if each of the members is infinite?
each member of an infinite set of sets of rationals?
each member of a denumerable set of sets if each of the members is denumerable?
each member of an infinite set of sets if each of the members is finite?
each member of an infinite set of finite sets of reals?
each member of an infinite set of sets of reals?
each member of an infinite set of two-element sets whose members are sets of reals?

### 2.3. Ordinal Numbers

Here we construct the ordinal numbers and extend the familiar ordering of the natural numbers. The construction is based on the notion of a well-ordered set.
2.38. Definition. Suppose $W$ is a well-ordered set with respect to the ordering $\leq$. We will use the notation $<$ in its familiar sense; we write $x<y$ to indicate that both $x \leq y$ and $x \neq y$. Also, in this case, we will agree to say that $x$ is less than $y$ and that $y$ is greater than $x$.

For $x \in W$ we define

$$
W(x)=\{y \in W: y<x\}
$$

and refer to $W(x)$ as the initial segment of $W$ determined by $x$.
The following is the Principle of Transfinite Induction.
2.39. Theorem. Let $W$ be a well-ordered set and let $S \subset W$ be defined as

$$
S:=\{x: W(x) \subset S \text { implies } x \in S\}
$$

Then $S=W$.
Proof. If $S \neq W$ then $W-S$ is a nonempty subset of $W$ and thus has a least element $x_{0}$. Then $W\left(x_{0}\right) \subset S$, which by hypothesis implies that $x_{0} \in S$ contradicting the fact that $x_{0} \in W-S$.

When applied to the well-ordered set $\mathbb{Z}$ of natural numbers, the hypothesis of Theorem 2.39 appears to differ in two ways from that of the Principle of Finite Induction, Theorem 2.1. First, it is not assumed that $1 \in S$ and second, in order to conclude that $x \in S$ we need to know that every predecessor of $x$ is in $S$ and not just its immediate predecessor. The first difference is illusory for suppose $a$ is the least element of $W$. Then $W(a)=\emptyset \subset S$ and thus $a \in S$. The second difference
is more significant because, unlike the case of $\mathbb{N}$, an element of an arbitrary wellordered set may not have an immediate predecessor.
2.40. Definition. A mapping $\varphi$ from a well-ordered set $V$ into a well-ordered set $W$ is order-preserving if $\varphi\left(v_{1}\right) \leq \varphi\left(v_{2}\right)$ whenever $v_{1}, v_{2} \in V$ and $v_{1} \leq v_{2}$. If, in addition, $\varphi$ is a bijection we will refer to it as an (order-preserving) isomorphism. Note that, in this case, $v_{1}<v_{2}$ implies $\varphi\left(v_{1}\right)<\varphi\left(v_{2}\right)$; in other words, an orderpreserving isomorphism is strictly order-preserving.

Note: We have slightly abused the notation by using the same symbol $\leq$ to indicate the ordering in both $V$ and $W$ above. But this should cause no confusion.
2.41. Lemma. If $\varphi$ is an order-preserving injection of a well-ordered set $W$ into itself, then

$$
w \leq \varphi(w)
$$

for each $w \in W$.
Proof. Set

$$
S=\{w \in W: \varphi(w)<w\}
$$

If $S$ is not empty, then it has a least element, say $a$. Thus $\varphi(a)<a$ and consequently $\varphi(\varphi(a))<\varphi(a)$ since $\varphi$ is an order-preserving injection; moreover, $\varphi(a) \notin S$ since $a$ is the least element of $S$. By the definition of $S$, this implies $\varphi(a) \leq \varphi(\varphi(a))$, which is a contradiction.
2.42. Corollary. If $V$ and $W$ are two well-ordered sets, then there is at most one isomorphism of $V$ onto $W$.

Proof. Suppose $f$ and $g$ are isomorphisms of $V$ onto $W$. Then $g^{-1} \circ f$ is an isomorphism of $V$ onto itself and hence $v \leq g^{-1} \circ f(v)$ for each $v \in V$. This implies that $g(v) \leq f(v)$ for each $v \in V$. Since the same argument is valid with the roles of $f$ and $g$ interchanged, we see that $f=g$.
2.43. Corollary. If $W$ is a well-ordered set, then $W$ is not isomorphic to an initial segment of itself

Proof. Suppose $a \in W$ and $W \xrightarrow{f} W(a)$ is an isomorphism. Since $w \leq$ $f(w)$ for each $w \in W$, in particular we have $a \leq f(a)$. Hence $f(a) \notin W(a)$, a contradiction.
2.44. Corollary. No two distinct initial segments of a well ordered set $W$ are isomorphic.

Proof. Since one of the initial segments must be an initial segment of the other, the conclusion follows from the previous result.
2.45. Definition. We define an ordinal number as an equivalence class of well-ordered sets with respect to order-preserving isomorphisms. If $W$ is a wellordered set, we denote the corresponding ordinal number by $\operatorname{ord}(W)$. We define a linear ordering on the class of ordinal numbers as follows: if $\mathfrak{v}=\operatorname{ord}(V)$ and $\mathfrak{w}=$ $\operatorname{ord}(W)$, then $\mathfrak{v}<\mathfrak{w}$ if and only if $V$ is isomorphic to an initial segment of $W$. The fact that this defines a linear ordering follows from the next result.
2.46. THEOREM. If $\mathfrak{v}$ and $\mathfrak{w}$ are ordinal numbers, then precisely one of the following holds:
(i) $\mathfrak{v}=\mathfrak{w}$
(ii) $\mathfrak{v}<\mathfrak{w}$
(iii) $\mathfrak{v}>\mathfrak{w}$

Proof. Let $V$ and $W$ be well-ordered sets representing $\mathfrak{v}$, $\mathfrak{w}$ respectively and let $\mathcal{F}$ denote the family of all order isomorphisms from an initial segment of V (or $V$ itself) onto either an initial segment of W (or $W$ itself). Recall that a mapping from a subset of $V$ into $W$ is a subset of $V \times W$. We may assume that $V \neq \emptyset \neq W$. If $v$ and $w$ are the least elements of $V$ and $W$ respectively, then $\{(v, w)\} \in \mathcal{F}$ and so $\mathcal{F}$ is not empty. Ordering $\mathcal{F}$ by inclusion, we see that any linearly ordered subset $S$ of $\mathcal{F}$ has an upper bound; indeed the union of the subsets of $V \times W$ corresponding to the elements of $S$ is easily seen to be an order isomorphism and thus an upper bound for $S$. Therefore we may employ Zorn's lemma to conclude that $\mathcal{F}$ has a maximal element, say $h$. Since $h \in \mathcal{F}$, it is an order isomorphism and $h \subset V \times W$. If domain $h$ and range $h$ were initial segments say $V_{x}$ and $W_{y}$ of $V$ and $W$, then $h^{*}:=h \cup\{(x, y)\}$ would contradict the maximality of $h$ unless domain $h=V$ or range $h=W$. If domain $h=V$, then either range $h=W$ (i.e., $\mathfrak{v}<\mathfrak{w}$ ) or range $h$ is an initial segment of $W$, (i.e., $\mathfrak{v}=\mathfrak{w}$ ). If domain $h \neq V$, then domain $h$ is an initial segment of $V$ and range $h=W$ and the existence of $h^{-1}$ in this case establishes $\mathfrak{v}>\mathfrak{w}$.

### 2.47. Theorem. The class of ordinal numbers is well-ordered.

Proof. Let $S$ be a nonempty set of ordinal numbers. Let $\alpha \in S$ and set

$$
T=\{\beta \in S: \beta<\alpha\}
$$

If $T=\emptyset$, then $\alpha$ is the least element of $S$. If $T \neq \emptyset$, let $W$ be a well-ordered set such that $\alpha=\operatorname{ord}(W)$. For each $\beta \in T$ there is a well-ordered set $W_{\beta}$ such that $\beta=$
$\operatorname{ord}\left(W_{\beta}\right)$, and there is a unique $x_{\beta} \in W$ such that $W_{\beta}$ is isomorphic to the initial segment $W\left(x_{\beta}\right)$ of $W$. The nonempty subset $\left\{x_{\beta}: \beta \in T\right\}$ of $W$ has a least element $x_{\beta_{0}}$. The element $\beta_{0} \in T$ is the least element of $T$ and thus the least element of $S$.

### 2.48. Corollary. The cardinal numbers are comparable.

Proof. Suppose $a$ is a cardinal number. Then, the set of all ordinals whose cardinal number is $a$ forms a well-ordered set that has a least element, call it $\alpha(a)$. The ordinal $\alpha(a)$ is called the initial ordinal of $a$. Suppose $b$ is another cardinal number and let $W(a)$ and $W(b)$ be the well-ordered sets whose ordinal numbers are $\alpha(a)$ and $\alpha(b)$, respectively. Either one of $W(a)$ or $W(b)$ is isomorphic to an initial segment of the other if $a$ and $b$ are not of the same cardinality. Thus, one of the sets $W(a)$ and $W(b)$ is equivalent to a subset of the other.
2.49. Corollary. Suppose $\alpha$ is an ordinal number. Then

$$
\alpha=\operatorname{ord}(\{\beta: \beta \text { is an ordinal number and } \beta<\alpha\})
$$

Proof. Let $W$ be a well-ordered set such that $\alpha=\operatorname{ord}(W)$. Let $\beta<\alpha$ and let $W(\beta)$ be the initial segment of $W$ whose ordinal number is $\beta$. It is easy to verify that this establishes an isomorphism between the elements of $W$ and the set of ordinals less than $\alpha$.

We may view the positive integers $\mathbb{N}$ as ordinal numbers in the following way. Set

$$
\begin{aligned}
1= & \operatorname{ord}(\{1\}) \\
2= & \operatorname{ord}(\{1,2\}) \\
3= & \operatorname{ord}(\{1,2,3\}) \\
& \vdots \\
\omega= & \operatorname{ord}(\mathbb{N})
\end{aligned}
$$

We see that

$$
\begin{equation*}
n<\omega \text { for each } n \in \mathbb{N} . \tag{2.7}
\end{equation*}
$$

If $\beta=\operatorname{ord}(W)<\omega$, then $W$ must be isomorphic to an initial segment of $\mathbb{N}$, i.e., $\beta=n$ for some $n \in \mathbb{N}$. Thus $\omega$ is the first ordinal number such that (2.7) holds and is thus the first infinite ordinal.

Consider the set of all ordinal numbers that have either finite or denumerable cardinal numbers and observe that this forms a well-ordered set. We denote the ordinal number of this set by $\Omega$. It can be shown that $\Omega$ is the first nondenumerable ordinal number, see Exercise 2 at the end of this section. The cardinal number of $\Omega$ is designated by $\aleph_{1}$. We have shown that $2^{\aleph_{0}}>\aleph_{0}$ and that $2^{\aleph_{0}}=c$. A fundamental question that remains open is whether $2^{\aleph_{0}}=\aleph_{1}$. The assertion that this equality holds is known as the continuum hypothesis. The work of Gödel [31] and Cohen [16], [17] shows that the continuum hypothesis and its negation are both consistent with the standard axioms of set theory.

At this point we acknowledge the inadequacy of the intuitive approach that we have taken to set theory. In the statement of Theorem 2.47 we were careful to refer to the class of ordinal numbers. This is because the ordinal numbers must not be a set! Suppose, for a moment, that the ordinal numbers form a set, say $\mathcal{O}$. Then according to Theorem $2.47, \mathcal{O}$ is a well-ordered set. Let $\sigma=\operatorname{ord}(\mathcal{O})$. Since $\sigma \in \mathcal{O}$ we must conclude that $\mathcal{O}$ is isomorphic to an initial segment of itself, contradicting Corollary 2.43. For an enlightening discussion of this situation see the book by P. R. Halmos [33].

## Exercises for Section 2.3

1. If $E$ is a set of ordinal numbers, prove that that there is an ordinal number $\alpha$ such that $\alpha>\beta$ for each $\beta \in E$.
2. Prove that $\Omega$ is the smallest nondenumerable ordinal.
3. Prove that the cardinality of all open sets in $\mathbb{R}^{n}$ is $c$.
4. Prove that the cardinality of all countable intersections of open sets in $\mathbb{R}^{n}$ is $c$.
5. Prove that the cardinality of all sequences of real numbers is $c$.
6. Prove that there are uncountably many subsets of an infinite set that are infinite.

## CHAPTER 3

## Elements of Topology

### 3.1. Topological Spaces

The purpose of this short chapter is to provide enough point set topology for the development of the subsequent material in real analysis. An indepth treatment is not intended. In this section, we begin with basic concepts and properties of topological spaces.

Here, instead of the word "set," the word "space" appears for the first time. Often the word "space" is used to designate a set that has been endowed with a special structure. For example a vector space is a set, such as $\mathbb{R}^{n}$, that has been endowed with an algebraic structure. Let us now turn to a short discussion of topological spaces.
3.1. Definition. The pair $(X, \mathcal{T})$ is called a topological space where $X$ is a nonempty set and $\mathcal{T}$ is a family of subsets of $X$ satisfying the following three conditions:
(i) The empty set $\emptyset$ and the whole space $X$ are elements of $\mathcal{T}$,
(ii) If $\mathcal{S}$ is an arbitrary subcollection of $\mathcal{T}$, then

$$
\bigcup\{U: U \in \mathcal{S}\} \in \mathcal{T}
$$

(iii) If $\mathcal{S}$ is any finite subcollection of $\mathcal{T}$, then

$$
\bigcap\{U: U \in \mathcal{S}\} \in \mathcal{T}
$$

The collection $\mathcal{T}$ is called a topology for the space $X$ and the elements of $\mathcal{T}$ are called the open sets of $X$. An open set containing a point $x \in X$ is called a neighborhood of $x$. The interior of an arbitrary set $A \subset X$ is the union of all open sets contained in $A$ and is denoted by $A^{\circ}$. Note that $A^{\circ}$ is an open set and that it is possible for some sets to have an empty interior. A set $A \subset X$ is called closed if $X \backslash A:=\tilde{A}$ is open. The closure of a set $A \subset X$, denoted by $\bar{A}$, is

$$
\bar{A}=X \cap\{x: U \cap A \neq \emptyset \text { for each open set } U \text { containing } x\}
$$

and the boundary of $A$ is $\partial A=\bar{A} \backslash A^{\circ}$. Note that $A \subset \bar{A}$.

These definitions are fundamental and will be used extensively throughout this text.
3.2. Definition. A point $x_{0}$ is called a limit point of a set $A \subset X$ provided $A \cap U$ contains a point of $A$ different from $x_{0}$ whenever $U$ is an open set containing $x_{0}$. The definition does not require $x_{0}$ to be an element of $A$. We will use the notation $A^{*}$ to denote the set of limit points of $A$.
3.3. Examples. (i) If $X$ is any set and $\mathcal{T}$ the family of all subsets of $X$, then $\mathcal{T}$ is called the discrete topology. It is the largest topology (in the sense of inclusion) that $X$ can possess. In this topology, all subsets of $X$ are open.
(ii) The indiscrete is where $\mathcal{T}$ is taken as only the empty set $\emptyset$ and $X$ itself; it is obviously the smallest topology on $X$. In this topology, the only open sets are $X$ and $\emptyset$.
(iii) Let $X=\mathbb{R}^{n}$ and let $\mathcal{T}$ consist of all sets $U$ satisfying the following property: for each point $x \in U$ there exists a number $r>0$ such that $B(x, r) \subset U$. Here, $B(x, r)$ denotes the ball or radius $r$ centered at $x$; that is,

$$
B(x, r)=\{y:|x-y|<r\}
$$

It is easy to verify that $\mathcal{T}$ is a topology. Note that $B(x, r)$ itself is an open set. This is true because if $y \in B(x, r)$ and $t=r-|y-x|$, then an application of the triangle inequality shows that $B(y, t) \subset B(x, r)$. Of course, for $n=1$, we have that $B(x, r)$ is an open interval in $\mathbb{R}$.
(iv) Let $X=[0,1] \cup(1,2)$ and let $\mathcal{T}$ consist of $\{0\}$ and $\{1\}$ along with all open sets (open relative to $\mathbb{R}$ ) in $(0,1) \cup(1,2)$. Then the open sets in this topology contain, in particular, $[0,1]$ and $[1,2)$.
3.4. Definition. Suppose $Y \subset X$ and $\mathcal{T}$ is a topology for $X$. Then it is easy to see that the family $\mathcal{S}$ of sets of the form $Y \cap U$ where $U$ ranges over all elements of $\mathcal{T}$ satisfies the conditions for a topology on $Y$. The topology formed in this way is called the induced topology or equivalently, the relative topology on $Y$. The space $Y$ is said to inherit the topology from its parent space $X$.
3.5. Example. Let $X=\mathbb{R}^{2}$ and let $\mathcal{T}$ be the topology described in (iii) above. Let $Y=\mathbb{R}^{2} \cap\left\{x=\left(x_{1}, x_{2}\right): x_{2} \geq 0\right\} \cup\left\{x=\left(x_{1}, x_{2}\right): x_{1}=0\right\}$. Thus, $Y$ is the upper half-space of $\mathbb{R}^{2}$ along with both the horizontal and vertical axes. All intervals $I$ of the form $I=\left\{x=\left(x_{1}, x_{2}\right): x_{1}=0, a<x_{2}<b<0\right\}$, where $a$ and $b$ are arbitrary negative real numbers, are open in the induced topology on $Y$, but none of them is open in the topology on $X$. However, all intervals $J$ of the form
$J=\left\{x=\left(x_{1}, x_{2}\right): x_{1}=0, a \leq x_{2} \leq b\right\}$ are closed both in the relative topology and the topology on $X$.
3.6. Theorem. Let $(X, \mathcal{T})$ be a topological space. Then
(i) The union of an arbitrary collection of open sets is open.
(ii) The intersection of a finite number of open sets is open.
(iii) The union of a finite number of closed sets is closed.
(iv) The intersection of an arbitrary collection of closed sets is closed.
(v) $\overline{A \cup B}=\bar{A} \cup \bar{B}$ whenever $A, B \subset X$.
(vi) If $\left\{A_{\alpha}\right\}$ is an arbitrary collection of subsets of $X$, then

$$
\bigcup_{\alpha} \overline{A_{\alpha}} \subset \overline{\bigcup_{\alpha} A_{\alpha}}
$$

(vii) $\overline{A \cap B} \subset \bar{A} \cap \bar{B}$ whenever $A, B \subset X$.
(viii) $A$ set $A \subset X$ is closed if and only if $A=\bar{A}$.
(ix) $\bar{A}=A \cup A^{*}$

Proof. Parts (i) and (ii) constitute a restatement of the definition of a topological space. Parts (iii) and (iv) follow from (i) and (ii) and de Morgan's laws, 1.5.
(v) Since $A \subset A \cup B$, we have $\bar{A} \subset \overline{A \cup B}$. Similarly, $\bar{B} \subset \overline{A \cup B}$, thus proving $\overline{A \cup B} \supset \bar{A} \cup \bar{B}$. By contradiction, suppose the converse if not true. Then there exists $x \in \overline{A \cup B}$ with $x \notin \bar{A} \cup \bar{B}$ and therefore there exist open sets $U$ and $V$ containing $x$ such that $U \cap A=\emptyset=V \cap B$. However, since $U \cap V$ is an open set containing $x$, it follows that

$$
\emptyset \neq(U \cap V) \cap(A \cup B) \subset(U \cap A) \cup(V \cap B)=\emptyset
$$

a contradiction.
(vi) This follows from the same reasoning used to establish the first part of (v).
(vii) This is immediate from definitions.
(viii) If $A=\bar{A}$, then $\tilde{A}$ is open (and thus $A$ is closed) because $x \notin \bar{A}$ implies that there exists an open set $U$ containing $x$ with $U \cap A=\emptyset$; that is, $U \subset \tilde{A}$. Conversely, if $A$ is closed and $x \in \tilde{A}$, then $x$ belongs to some open set $U$ with $U \subset \tilde{A}$. Thus, $U \cap A=\emptyset$ and therefore $x \notin \bar{A}$. This proves $\tilde{A} \subset(\bar{A})^{\sim}$ or $\bar{A} \subset A$. But always $A \subset \bar{A}$ and hence, $A=\bar{A}$.
(ix) is left as Exercise 2, Section 3.1.
3.7. Definition. Let $(X, \mathcal{T})$ be a topological space and $\left\{x_{i}\right\}_{i=1}^{\infty}$ a sequence in $X$. The sequence is said to converge to $x_{0} \in X$ if for each neighborhood $U$ of $x_{0}$ there is a positive integer $N$ such that $x_{i} \in U$ whenever $i \geq N$.

It is important to observe that the structure of a topological space is so general that a sequence could possibly have more than one limit. For example, every sequence in the space with the indiscrete topology (Example 3.3 (ii)) converges to every point in $X$. This cannot happen if an additional restriction is placed on the topological structure, as in the following definition. (Also note that the only sequences that converge in the discrete topology are those that are eventually constant.)
3.8. Definition. A topological space $X$ is said to be a Hausdorff space if for each pair of distinct points $x_{1}, x_{2} \in X$ there exist disjoint open sets $U_{1}$ and $U_{2}$ containing $x_{1}$ and $x_{2}$ respectively. That is, two distinct points can be separated by disjoint open sets.
3.9. Definition. Suppose $(X, \mathcal{T})$ and $(Y, \mathcal{S})$ are topological spaces. A function $f: X \rightarrow Y$ is said to be continuous at $x_{0} \in X$ if for each neighborhood $V$ containing $f\left(x_{0}\right)$ there is a neighborhood $U$ of $x_{0}$ such that $f(U) \subset V$. The function $f$ is said to be continuous on $X$ if it is continuous at each point $x \in X$.

The proof of the next result is given as Exercise 4, Section 3.1.
3.10. Theorem. Let $(X, \mathcal{T})$ and $(Y, \mathcal{S})$ be topological spaces. Then for a function $f: X \rightarrow Y$, the following statements are equivalent:
(i) $f$ is continuous.
(ii) $f^{-1}(V)$ is open in $X$ for each open set $V$ in $Y$.
(iii) $f^{-1}(K)$ is closed in $X$ for each closed set $K$ in $Y$.
3.11. Definition. A collection of open sets, $\mathcal{F}$, in a topological space $X$ is said to be an open cover of a set $A \subset X$ if

$$
A \subset \bigcup_{U \in \mathcal{F}} U
$$

The family $\mathcal{F}$ is said to admit a subcover, $\mathcal{G}$, of $A$ if $\mathcal{G} \subset \mathcal{F}$ and $\mathcal{G}$ is a cover of A. A subset $K \subset X$ is called compact if each open cover of $K$ possesses a finite subcover of $K$. A space $X$ is said to be locally compact if each point of $X$ is contained in some open set whose closure is compact.

It is easy to give illustrations of sets that are not compact. For example, it is readily seen that the set $A=(0,1]$ in $\mathbb{R}$ is not compact since the collection of open intervals of the form $(1 / i, 2), i=1,2, \ldots$, provides an open cover of $A$ admits no finite subcover. On the other hand, it is true that $[0,1]$ is compact, but the proof is not obvious. The reason for this is that the definition of compactness is usually
not easy to employ directly. Later, in the context of metric spaces (Section 3.3), we will find other ways of dealing with compactness.

The following two propositions reveal some basic connections between closed and compact subsets.
3.12. Proposition. Let $(X, \mathcal{T})$ be a topological space. If $A$ and $K$ are respectively closed and compact subsets of $X$ with $A \subset K$, then $A$ is compact.

Proof. If $\mathcal{F}$ is an open cover of $A$, then the elements of $\mathcal{F}$ along with $X \backslash A$ form an open cover of $K$. This open cover has a finite subcover, $\mathcal{G}$, of $K$ since $K$ is compact. The set $X \backslash A$ may possibly be an element of $\mathcal{G}$. If $X \backslash A$ is not a member of $\mathcal{G}$, then $\mathcal{G}$ is a finite subcover of $A$; if $X \backslash A$ is a member of $\mathcal{G}$, then $\mathcal{G}$ with $X \backslash A$ omitted is a finite subcover of $A$.
3.13. Proposition. A compact subset of a Hausdorff space $(X, \mathcal{T})$ is closed.

Proof. We will show that $X \backslash K$ is open where $K \subset X$ is compact. Choose a fixed $x_{0} \in X \backslash K$ and for each $y \in K$, let $V_{y}$ and $U_{y}$ denote disjoint neighborhoods of $y$ and $x_{0}$ respectively. The family

$$
\mathcal{F}=\left\{V_{y}: y \in K\right\}
$$

forms an open cover of $K$. Hence, $\mathcal{F}$ possesses a finite subcover, say $\left\{V_{y_{i}}: i=\right.$ $1,2, \ldots, N\}$. Since $V_{y_{i}} \cap U_{y_{i}}=\emptyset, i=1,2, \ldots, N$, it follows that $\bigcap_{i=1}^{N} U_{y_{i}} \cap \bigcup_{i=1}^{N} V_{y_{i}}=\emptyset$. Since $K \subset \cup_{i=1}^{N} V_{y_{i}}$ it follows that $\cap_{i=1}^{N} V_{y_{i}}$ is an open set containing $x_{0}$ that does not intersect $K$. Thus, $X \backslash K$ is an open set, as desired.

The characteristic property of a Hausdorff space is that two distinct points can be separated by disjoint open sets. The next result shows that a stronger property holds, namely, that a compact set and a point not in this compact set can be separated by disjoint open sets.
3.14. Proposition. Suppose $K$ is a compact subset of a Hausdorff space $X$ space and assume $x_{0} \notin K$. Then there exist disjoint open sets $U$ and $V$ containing $x_{0}$ and $K$ respectively.

Proof. This follows immediately from the preceding proof by taking

$$
U=\bigcap_{i=1}^{N} U_{y_{i}} \quad \text { and } \quad V=\bigcup_{i=1}^{N} V_{y_{i}} .
$$

3.15. Definition. A family $\left\{E_{\alpha}: \alpha \in I\right\}$ of subsets of a set $X$ is said to have the finite intersection property if for each finite subset $F \subset I$

$$
\bigcap_{\alpha \in F} E_{\alpha} \neq \emptyset
$$

3.16. Lemma. A topological space $X$ is compact if and only if every family of closed subsets of $X$ having the finite intersection property has a nonempty intersection.

Proof. First assume that $X$ is compact and let $\left\{C_{\alpha}\right\}$ be a family of closed sets with the finite intersection property. Then $\left\{U_{\alpha}\right\}:=\left\{X \backslash C_{\alpha}\right\}$ is a family, $\mathcal{F}$, of open sets. If $\bigcap_{\alpha} C_{\alpha}$ were empty, then $\mathcal{F}$ would form an open covering of $X$ and therefore the compactness of $X$ would imply that $\mathcal{F}$ has a finite subcover. This would imply that $\left\{C_{\alpha}\right\}$ has a finite subfamily with an empty intersection, contradicting the fact that $\left\{C_{\alpha}\right\}$ has the finite intersection property.

For the converse, let $\left\{U_{\alpha}\right\}$ be an open covering of $X$ and let $\left\{C_{\alpha}\right\}:=\left\{X \backslash U_{\alpha}\right\}$. If $\left\{U_{\alpha}\right\}$ had no finite subcover of $X$, then $\left\{C_{\alpha}\right\}$ would have the finite intersection property, and therefore, $\bigcap_{\alpha} C_{\alpha}$ would be nonempty, thus contradicting the assumption that $\left\{U_{\alpha\}}\right.$ is a covering of $X$.
3.17. Remark. An equivalent way of stating the previous result is as follows: A topological space $X$ is compact if and only if every family of closed subsets of $X$ whose intersection is empty has a finite subfamily whose intersection is also empty.
3.18. Theorem. Suppose $K \subset U$ are respectively compact and open sets in a locally compact Hausdorff space $X$. Then there is an open set $V$ whose closure is compact such that

$$
K \subset V \subset \bar{V} \subset U
$$

Proof. Since each point of $K$ is contained in an open set whose closure is compact, and since $K$ can be covered by finitely many such open sets, it follows that the union of these open sets, call it $G$, is an open set containing $K$ with compact closure. Thus, if $U=X$, the proof is compete.

Now consider the case $U \neq X$. Proposition 3.14 states that for each $x \in \widetilde{U}$ there is an open set $V_{x}$ such that $K \subset V_{x}$ and $x \notin \bar{V}_{x}$. Let $\mathcal{F}$ be the family of compact sets defined by

$$
\mathcal{F}:=\left\{\widetilde{U} \cap \bar{G} \cap \bar{V}_{x}: x \in \widetilde{U}\right\} .
$$

and observe that the intersection of all sets in $\mathcal{F}$ is empty, for otherwise, we would be faced with impossibility of some $x_{0} \in \widetilde{U} \cap \bar{G}$ that also belongs to $\bar{V}_{x_{0}}$. Lemma 3.16 (or Remark 3.17) implies there is some finite subfamily of $\mathcal{F}$ that has an empty intersection. That is, there exist points $x_{1}, x_{2}, \ldots, x_{k} \in \widetilde{U}$ such that

$$
\widetilde{U} \cap \bar{G} \cap \bar{V}_{x_{1}} \cap \cdots \cap \bar{V}_{x_{k}}=\emptyset .
$$

The set

$$
V=G \cap V_{x_{1}} \cap \cdots \cap V_{x_{k}}
$$

satisfies the conclusion of our theorem since

$$
K \subset V \subset \bar{V} \subset \bar{G} \cap \bar{V}_{x_{1}} \cap \cdots \cap \bar{V}_{x_{k}} \subset U
$$

## Exercises for Section 3.1

1. In a topological space $(X, \mathcal{T})$, prove that $\bar{A}=\overline{\bar{A}}$ whenever $A \subset X$.
2. Prove (ix) of Theorem 3.6.
3. Prove that $A^{*}$ is a closed set.
4. Prove Theorem 3.10.

### 3.2. Bases for a Topology

Often a topology is described in terms of a primitive family of sets, called a basis. We will give a brief description of this concept.
3.19. Definition. A collection $\mathcal{B}$ of open sets in a topological space $(X, \mathcal{T})$ is called a basis for the topology $\mathcal{T}$ if and only if $\mathcal{B}$ is a subfamily of $\mathcal{T}$ with the property that for each $U \in \mathcal{T}$ and each $x \in U$, there exists $B \in \mathcal{B}$ such that $x \in B \subset U$. A collection $\mathcal{B}$ of open sets containing a point $x$ is said to be a basis at $x$ if for each open set $U$ containing $x$ there is a $B \in \mathcal{B}$ such that $x \in B \subset U$. Observe that a collection $\mathcal{B}$ forms a basis for a topology if and only if it contains a basis at each point $x \in X$. For example, the collection of all sets $B(x, r), r>0, x \in \mathbb{R}^{n}$, provides a basis for the topology on $\mathbb{R}^{n}$ as described in (iii) of Example 3.3.

The following is a useful tool for generating a topology on a space $X$.
3.20. Proposition. Let $X$ be an arbitrary space. A collection $\mathcal{B}$ of subsets of $X$ is a basis for some topology on $X$ if and only if each $x \in X$ is contained in some $B \in \mathcal{B}$ and if $x \in B_{1} \cap B_{2}$, then there exists $B_{3} \in \mathcal{B}$ such that $x \in B_{3} \subset B_{1} \cap B_{2}$.

Proof. It is easy to verify that the conditions specified in the Proposition are necessary. To show that they are sufficient, let $\mathcal{T}$ be the collection of sets $U$ with the property that for each $x \in U$, there exists $B \in \mathcal{B}$ such that $x \in B \subset U$. It is easy to verify that $\mathcal{T}$ is closed under arbitrary unions. To show that it is closed under finite intersections, it is sufficient to consider the case of two sets. Thus, suppose $x \in U_{1} \cap U_{2}$, where $U_{1}$ and $U_{2}$ are elements of $\mathcal{T}$. There exist $B_{1}, B_{2} \in \mathcal{B}$ such that $x \in B_{1} \subset U_{1}$ and $x \in B_{2} \subset U_{2}$. We are given that there is $B_{3} \in \mathcal{B}$ such that $x \in B_{3} \subset B_{1} \cap B_{2}$, thus showing that $U_{1} \cap U_{2} \in \mathcal{T}$.
3.21. Definition. A topological space $(X, \mathcal{T})$ is said to satisfy the first axiom of countability if each point $x \in X$ has a countable basis $\mathcal{B}_{x}$ at $x$. It is said to satisfy the second axiom of countability if there is a countable basis $\mathcal{B}$.

The second axiom of countability obviously implies the first axiom of countability. The usual topology on $\mathbb{R}^{n}$, for example, satisfies the second axiom of countability.
3.22. Definition. A family $S$ of subsets of a topological space $(X, \mathcal{T})$ is called a subbase for the topology $\mathcal{T}$ if the family consisting of all finite intersections of members of $S$ forms a base for the topology $\mathcal{T}$.

In view of Proposition 3.20, every nonempty family of subsets of $X$ is the subbase for some topology on $X$. This leads to the concept of the product topology.
3.23. Definition. Given an index set $A$, consider the Cartesian product
$\prod_{\alpha \in A} X_{\alpha}$ where each $\left(X_{\alpha}, \mathcal{T}_{\alpha}\right)$ is a topological space. For each $\beta \in A$ there is a natural projection

$$
P_{\beta}: \prod_{\alpha \in A} X_{\alpha} \rightarrow X_{\beta}
$$

defined by $P_{\beta}(x)=x_{\beta}$ where $x_{\beta}$ is the $\beta^{t h}$ coordinate of $x$, (See (1.13) and its following remarks.) Consider the collection $S$ of subsets of $\prod_{\alpha \in A} X_{\alpha}$ given by

$$
P_{\alpha}^{-1}\left(V_{\alpha}\right)
$$

where $V_{\alpha} \in \mathcal{T}_{\alpha}$ and $\alpha \in A$. The topology formed by the subbase $S$ is called the product topology on $\prod_{\alpha \in A} X_{\alpha}$. In this topology, the projection maps $P_{\beta}$ are continuous.

It is easily seen that a function $f$ from a topological space $(Y, \mathcal{T})$ into a product space $\prod_{\alpha \in A} X_{\alpha}$ is continuous if and only if $\left(P_{\alpha} \circ f\right)$ is continuous for each $\alpha \in A$. Moreover, a sequence $\left\{x_{i}\right\}_{i=1}^{\infty}$ in a product space $\prod_{\alpha \in A} X_{\alpha}$ converges to a point $x_{0}$ of the product space if and only if the sequence $\left\{P_{\alpha}\left(x_{i}\right)\right\}_{i=1}^{\infty}$ converges to $P_{\alpha}\left(x_{0}\right)$ for each $\alpha \in A$. See Exercises 4 and 5 at the end of this section.

## Exercises for Section 3.2

1. Prove that the product topology on $\mathbb{R}^{n}$ agrees with the Euclidean topology on $\mathbb{R}^{n}$.
2. Suppose that $X_{i}, i=1,2$ satisfy the second axiom of countability. Prove that the product space $X_{1} \times X_{2}$ also satisfies the second axiom of countability.
3. Let $(X, \mathcal{T})$ be a topological space and let $f: X \rightarrow \mathbb{R}$ and $g: X \rightarrow \mathbb{R}$ be continuous functions. Define $F: X \rightarrow \mathbb{R} \times \mathbb{R}$ by

$$
F(x)=(f(x), g(x)), \quad x \in X
$$

Prove that $F$ is continuous.
4. Show that a function $f$ from a topological space $(X, \mathcal{T})$ into a product space $\prod_{\alpha \in A} X_{\alpha}$ is continuous if and only if $\left(P_{\alpha} \circ f\right)$ is continuous for each $\alpha \in A$.
5. Prove that a sequence $\left\{x_{i}\right\}_{i=1}^{\infty}$ in a product space $\prod_{\alpha \in A} X_{\alpha}$ converges to a point $x_{0}$ of the product space if and only if the sequence $\left\{P_{\alpha}\left(x_{i}\right)\right\}_{i=1}^{\infty}$ converges to $P_{\alpha}\left(x_{0}\right)$ for each $\alpha \in A$.

### 3.3. Metric Spaces

Metric spaces are used extensively throughout analysis. The main purpose of this section is to introduce basic definitions.

We already have mentioned two structures placed on sets that deserve the designation space, namely, vector space and topological space. We now come to our third structure, that of a metric space.
3.24. Definition. A metric space is an arbitrary set $X$ endowed with a metric $\rho: X \times X \rightarrow[0, \infty)$ that satisfies the following properties for all $x, y$ and $z$ in $X$ :
(i) $\rho(x, y)=0$ if and only if $x=y$,
(ii) $\rho(x, y)=\rho(y, x)$,
(iii) $\rho(x, y) \leq \rho(x, z)+\rho(z, y)$.

We will write $(X, \rho)$ to denote the metric space $X$ endowed with a metric $\rho$. Often the metric $\rho$ is called the distance function and a reasonable name for property (iii) is the triangle inequality. If $Y \subset X$, then the metric space $(Y, \rho L(Y \times Y))$ is called the subspace induced by $(X, \rho)$.

The following are easily seen to be metric spaces.

### 3.25. Example.

(i) Let $X=\mathbb{R}^{n}$ and with $x=\left(x_{1}, \ldots, x_{n}\right), y=\left(y_{1}, \ldots, y_{n}\right) \in \mathbb{R}^{n}$, define

$$
\rho(x, y)=\left(\sum_{i=1}^{n}\left|x_{i}-y_{i}\right|^{2}\right)^{1 / 2}
$$

(ii) Let $X=\mathbb{R}^{n}$ and with $x=\left(x_{1}, \ldots, x_{n}\right), y=\left(y_{1}, \ldots, y_{n}\right) \in \mathbb{R}^{n}$, define

$$
\rho(x, y)=\max \left\{\left|x_{i}-y_{i}\right|: i=1,2, \ldots, n\right\} .
$$

(iii) The discrete metric on and arbitrary set $X$ is defined as follows: for $x, y \in$ $X$,

$$
\rho(x, y)= \begin{cases}1 & \text { if } x \neq y \\ 0 & \text { if } x=y\end{cases}
$$

(iv) Let $X$ denote the space of all continuous functions defined on $[0,1]$ and for $f, g \in C(X)$, let

$$
\rho(f, g)=\int_{0}^{1}|f(t)-g(t)| d t
$$

(v) Let $X$ denote the space of all continuous functions defined on $[0,1]$ and for $f, g \in C(X)$, let

$$
\rho(f, g)=\max \{|f(x)-g(x)|: x \in[0,1]\}
$$

3.26. Definition. If $X$ is a metric space with metric $\rho$, the open ball centered at $x \in X$ with radius $r>0$ is defined as

$$
B(x, r)=X \cap\{y: \rho(x, y)<r\}
$$

The closed ball is defined as

$$
\bar{B}(x, r):=X \cap\{y: \rho(x, y) \leq r\} .
$$

In view of the triangle inequality, the family $S=\{B(x, r): x \in X, r>0\}$ forms a basis for a topology $\mathcal{T}$ on $X$ called the topology induced by $\rho$. The two metrics in $\mathbb{R}^{n}$ defined in Examples 3.25, (i) and (ii) induce the same topology on $\mathbb{R}^{n}$. Two metrics on a set $X$ are said to be topologically equivalent if they induce the same topology on $X$.
3.27. Definition. Using the notion of convergence given in Definition 3.7, p.39, the reader can easily verify that the convergence of a sequence $\left\{x_{i}\right\}_{i=1}^{\infty}$ in a metric space $(X, \rho)$ becomes the following:

$$
\lim _{i \rightarrow \infty} x_{i}=x_{0}
$$

if and only if for each positive number $\varepsilon$ there is a positive integer $N$ such that

$$
\rho\left(x_{i}, x_{0}\right)<\varepsilon \quad \text { whenever } \quad i \geq N
$$

We often write $x_{i} \rightarrow x_{0}$ for $\lim _{i \rightarrow \infty} x_{i}=x_{0}$.
The notion of a fundamental or a Cauchy sequence is not a topological one and requires a separate definition:
3.28. Definition. A sequence $\left\{x_{i}\right\}_{i=1}^{\infty}$ is called Cauchy if for every $\varepsilon>0$, there exists a positive integer $N$ such that $\rho\left(x_{i}, x_{j}\right)<\varepsilon$ whenever $i, j \geq N$. The notation for this is

$$
\lim _{i, j \rightarrow \infty} \rho\left(x_{i}, x_{j}\right)=0
$$

Recall the definition of continuity given in Definition 3.9. In a metric space, it is convenient to have the following characterization whose proof is left as an exercise.
3.29. Theorem. If $(X, \rho)$ and $(Y, \sigma)$ are metric spaces, then a mapping $f: X \rightarrow Y$ is continuous at $x \in X$ if for each $\varepsilon>0$, there exists $\delta>0$ such that $\sigma[f(x), f(y)]<\varepsilon$ whenever $\rho(x, y)<\delta$.
3.30. Definition. If $X$ and $Y$ are topological spaces and if $f: X \rightarrow Y$ is a bijection with the property that both $f$ and $f^{-1}$ are continuous, then $f$ is called a homeomorphism and the spaces $X$ and $Y$ are said to be homeomorphic. A substantial part of topology is devoted to the investigation of those properties that remain unchanged under the action of a homeomorphism. For example, in view of Exercise 12 at the end of this section, it follows that if $U \subset X$ is open, then so is $f(U)$ whenever $f: X \rightarrow Y$ is a homeomorphism; that is, the property of being open is a topological invariant. Consequently, so is closedness. But of course, not all properties are topological invariants. For example, the distance between two points might be changed under a homeomorphism. A mapping that preserves distances, that is, one for which

$$
\sigma[f(x), f(y)]=\rho(x, y)
$$

for all $x, y \in X$ is called an isometry. In particular, it is a homeomorphism, The spaces $X$ and $Y$ are called isometric if there exists a surjection $f: X \rightarrow Y$ that is an isometry. In the context of metric space topology, isometric spaces can be regarded as being identical.

It is easy to verify that a convergent sequence in a metric space is Cauchy, but the converse need not be true. For example, the metric space, $\mathbb{Q}$, consisting of the rational numbers endowed with the usual metric on $\mathbb{R}$, possesses Cauchy sequences that do not converge to elements in $\mathbb{Q}$. If a metric space has the property that every Cauchy sequence converges (to an element of the space), the space is said to be complete. Thus, the metric space of rational numbers is not complete, whereas the real numbers are complete. However, we can apply the technique that was employed in the construction of the real numbers (see Section 2.1, p.13) to complete an arbitrary metric space. A precise statement of this is incorporated in the following theorem, whose proof is left as Exercise 2, Section 3.4.
3.31. Theorem. If $(X, \rho)$ is a metric space, there exists a complete metric space $\left(X^{*}, \rho^{*}\right)$ in which $X$ is isometrically embedded as a dense subset.

In the statement, the notion of a dense set is used. This notion is a topological one. In a topological space $(X, \mathcal{T})$, a subset $A$ of $X$ is said to be a dense subset of $X$ if $X=\bar{A}$.

## Exercises for Section 3.3

1. In a metric space, prove that $B(x, \rho)$ is an open set and that $\bar{B}(x, \rho)$ is closed. Is $\bar{B}(x, \rho)=\overline{B(x, \rho)}$ ?
2. Suppose $X$ is a complete metric space. Show that if $F_{1} \supset F_{2} \supset \ldots$ are nonempty closed subsets of $X$ with diameter $F_{i} \rightarrow 0$, then there exists $x \in X$ such that

$$
\bigcap_{i=1}^{\infty} F_{i}=\{x\}
$$

3. Suppose $(X, \rho)$ and $(Y, \sigma)$ are metric spaces with $X$ compact and $Y$ complete. Let $C(X, Y)$ denote the space of all continuous mappings $f: X \rightarrow Y$. Define a metric on $C(X, Y)$ by

$$
d(f, g)=\sup \{\sigma(f(x), g(x)): x \in X\} .
$$

Prove that $C(X, Y)$ is a complete metric space.
4. Let ( $X_{1}, \rho_{1}$ ) and ( $X_{2}, \rho_{2}$ ) be metric spaces and define metrics on $X_{1} \times X_{2}$ as follows: For $\boldsymbol{x}=\left(x_{1}, x_{2}\right), \boldsymbol{y}:=\left(y_{1}, y_{2}\right) \in X_{1} \times X_{2}$, let
(a) $d_{1}(\boldsymbol{x}, \boldsymbol{y}):=\rho_{1}\left(x_{1}, y_{1}\right)+\rho_{2}\left(x_{2}, y_{2}\right)$
(b) $d_{2}(\boldsymbol{x}, \boldsymbol{y}):=\sqrt{\left(\rho_{1}\left(x_{1}, y_{1}\right)\right)^{2}+\left(\rho_{2}\left(x_{2}, y_{2}\right)\right)^{2}}$
(i) Prove that $d_{1}$ and $d_{2}$ define identical topologies.
(ii) Prove that $\left(X_{1} \times X_{2}, d_{1}\right)$ is complete if and only if $X_{1}$ and $X_{2}$ are complete.
(iii) Prove that ( $X_{1} \times X_{2}, d_{1}$ ) is compact if and only if $X_{1}$ and $X_{2}$ are compact.
5. Suppose $A$ is a subset of a metric space $X$. Prove that a point $x_{0} \notin A$ is a limit point of $A$ if and only if there is a sequence $\left\{x_{i}\right\}$ in $A$ such that $x_{i} \rightarrow x_{0}$.
6. Prove that a closed subset of a complete metric space is a complete metric space.
7. A mapping $f: X \rightarrow X$ with the property that there exists a number $0<K<1$ such that $\rho(f(x), f(y))<K \rho(x, y)$ for all $x \neq y$ is called a contraction. Prove that a contraction on a complete metric space has a unique fixed point.
8. Suppose $(X, \rho)$ is metric space and consider a mapping from $X$ into itself, $f: X \rightarrow X$. A point $x_{0} \in X$ is called a fixed point for $f$ if $f\left(x_{0}\right)=x_{0}$. Prove that if $X$ is compact and $f$ has the property that $\rho(f(x), f(y))<\rho(x, y)$ for all $x \neq y$, then $f$ has a unique fixed point.
9. As on p.294, a mapping $f: X \rightarrow X$ with the property that $\rho(f(x), f(y))=\rho(x, y)$ for all $x, y \in X$ is called an isometry. If $X$ is compact, prove that an isometry is a surjection. Is compactness necessary?
10. Show that a metric space $X$ is compact if and only if every continuous real-valued function on $X$ attains a maximum value.
11. If $X$ and $Y$ are topological spaces, prove $f: X \rightarrow Y$ is continuous if and only if $f^{-1}(U)$ is open whenever $U \subset Y$ is open.
12. Suppose $f: X \rightarrow Y$ is surjective and a homeomorphism. Prove that if $U \subset X$ is open, then so is $f(U)$.
13. If $X$ and $Y$ are topological spaces, show that if $f: X \rightarrow Y$ is continuous, then $f\left(x_{i}\right) \rightarrow f\left(x_{0}\right)$ whenever $\left\{x_{i}\right\}$ is a sequence that converges to $x_{0}$. Show that the converse is true if $X$ and $Y$ are metric spaces.
14. Prove that a subset $C$ of a metric space $X$ is closed if and only if every convergent sequence $\left\{x_{i}\right\}$ in $C$ converges to a point in $C$.
15. Prove that $C[0,1]$ is not a complete space when endowed with the metric given in (iv) of Example 3.25, p. 45.
16. Prove that in a topological space $(X, \mathcal{T})$, if $A$ is dense in $B$ and $B$ is dense in $C$, then $A$ is dense in $C$.
17. A metric space is said to be separable if it has a countable dense subset.
(i) Show that $\mathbb{R}^{n}$ with its usual topology is separable.
(ii) Prove that a metric space is separable if and only if it satisfies the second axiom of countability.
(iii) Prove that a subspace of a separable metric space is separable.
(iv) Prove that if a metric space $X$ is separable, then $\operatorname{card} X \leq c$.
18. Let $(X, \varrho)$ be a metric space, $Y \subset X$, and let $(Y, \varrho\llcorner(Y \times Y))$ be the induced subspace. Prove that if $E \subset Y$, then the closure of $E$ in the subspace $Y$ is the same as the closure of $E$ in the space $X$ intersected with $Y$.
19. Prove that the discrete metric on $X$ induces the discrete topology on $X$.

### 3.4. Meager Sets in Topology


#### Abstract

Throughout this book, we will encounter several ways of describing the "size" of a set. In Chapter 2 the size of a set was described in terms of its cardinality. Later, we will discuss other methods. The notion of a nowhere dense set and its related concept, that of a set being of the first category, are ways of saying that a set is "meager" in the topological sense. In this section we shall prove one of the main results involving these concepts, the Baire Category Theorem, which asserts that a complete metric space is not meager.


Recall Definition 3.24 in which a subset $S$ of a metric space $(X, \rho)$ is endowed with the induced topology. The metric placed on $S$ is obtained by restricting the metric $\rho$ to $S \times S$. Thus, the distance between any two points $x, y \in S$ is defined as $\rho(x, y)$, which is the distance between $x, y$ as points of $X$.

As a result of the definition, a subset $U \subset S$ is open in $S$ if for each $x \in U$, there exists $r>0$ such that if $y \in S$ and $\rho(x, y)<r$, then $y \in U$. In other words, $B(x, r) \cap S \subset U$ where $B(x, r)$ is taken as the ball in $X$. Thus, it is easy to see that $U$ is open in $S$ if and only if there exists an open set $V$ in $X$ such that $U=V \cap S$. Consequently, a set $F \subset S$ is closed relative to $S$ if and only if $F=C \cap S$ for some
closed set $C$ in $S$. Moreover, the closure of a set $E$ relative to $S$ is $\bar{E} \cap S$, where $\bar{E}$ denotes the closure of $E$ in $X$. This is true because if a point $x$ is in the closure of $E$ in $X$, then it is a point in the closure of $E$ in $S$ if it belongs to $S$.
3.32. Definitions. A subset $E$ of a metric space $X$ is said to be dense in an open set $U$ if $\bar{E} \supset U$. Also, a set $E$ is defined to be nowhere dense if it is not dense in any open subset $U$ of $X$. Alternatively, we could say that $E$ is nowhere dense if $\bar{E}$ does not contain any open set. For example, the integers comprise a nowhere dense set in $\mathbb{R}$, whereas the set $Q \cap[0,1]$ is not nowhere dense in $\mathbb{R}$. A set $E$ is said to be of first category in $X$ if it is the union of a countable collection of nowhere dense sets. A set that is not of the first category is said to be of the second category in $X$.

We now proceed to investigate a fundamental result related to these concepts.
3.33. Theorem (Baire Category Theorem). A complete metric space $X$ is not the union of a countable collection of nowhere dense sets. That is, a complete metric space is of the second category.

Before going on, it is important to examine the statement of the theorem in various contexts. For example, let $X$ be the integers endowed with the metric induced from $\mathbb{R}$. Thus, $X$ is a complete metric space and therefore, by the Baire Category Theorem, it is of the second category. At first, this may seem counter intuitive, since $X$ is the union of a countable collection of points. But remember that a point in this space is an open set, and therefore is not nowhere dense. However, if $X$ is viewed as a subset of $\mathbb{R}$ and not as a space in itself, then indeed, $X$ is the union of a countable number of nowhere dense sets.

Proof. Assume by contradiction, that $X$ is of the first category. Then there exists a countable collection of nowhere dense sets $\left\{E_{i}\right\}$ such that

$$
X=\bigcup_{i=1}^{\infty} E_{i}
$$

Let $B\left(x_{1}, r_{1}\right)$ be an open ball with radius $r_{1}<1$. Since $E_{1}$ is not dense in any open set, it follows that $B\left(x_{1}, r_{1}\right) \backslash \overline{E_{1}} \neq \emptyset$. This is a nonempty open set, and therefore there exists a ball $B\left(x_{2}, r_{2}\right) \subset B\left(x_{1}, r_{1}\right) \backslash \overline{E_{1}}$ with $r_{2}<\frac{1}{2} r_{1}$. In fact, by also choosing $r_{2}$ smaller than $r_{1}-\rho\left(x_{1}, x_{2}\right)$, we may assume that $\bar{B}\left(x_{2}, r_{2}\right) \subset B\left(x_{1}, r_{1}\right) \backslash \overline{E_{1}}$. Similarly, since $E_{2}$ is not dense in any open set, we have that $B\left(x_{2}, r_{2}\right) \backslash \overline{E_{2}}$ is a nonempty open set. As before, we can find a closed ball with center $x_{3}$ and radius

$$
r_{3}<\frac{1}{2} r_{2}<\frac{1}{2^{2}} r_{1}
$$

$$
\begin{aligned}
\bar{B}\left(x_{3}, r_{3}\right) & \subset B\left(x_{2}, r_{2}\right) \backslash \overline{E_{2}} \\
& \subset\left(B\left(x_{1}, r_{1}\right) \backslash \overline{E_{1}}\right) \backslash \overline{E_{2}} \\
& =B\left(x_{1}, r_{1}\right) \backslash \bigcup_{j=1}^{2} \overline{E_{j}} .
\end{aligned}
$$

Proceeding inductively, we obtain a nested sequence $\bar{B}\left(x_{1}, r_{1}\right) \supset B\left(x_{1}, r_{1}\right) \supset$ $\bar{B}\left(x_{2}, r_{2}\right) \supset B\left(x_{2}, r_{2}\right) \ldots$ with $r_{i}<\frac{1}{2^{i}} r_{1} \rightarrow 0$ such that

$$
\begin{equation*}
\bar{B}\left(x_{i+1}, r_{i+1}\right) \subset B\left(x_{1}, r_{1}\right) \backslash \bigcup_{j=1}^{i} \overline{E_{j}} \tag{3.1}
\end{equation*}
$$

for each $i$. Now, for $i, j>N$, we have $x_{i}, x_{j} \in B\left(x_{N}, r_{N}\right)$ and therefore $\rho\left(x_{i}, x_{j}\right) \leq$ $2 r_{N}$. Thus, the sequence $\left\{x_{i}\right\}$ is Cauchy in $X$. Since $X$ is assumed to be complete, it follows that $x_{i} \rightarrow x$ for some $x \in X$. For each positive integer $N, x_{i} \in B\left(x_{N}, r_{N}\right)$ for $i \geq N$. Hence, $x \in \bar{B}\left(x_{N}, r_{N}\right)$ for each positive integer $N$. For each positive integer $i$ it follows from (3.1) that

$$
x \in \bar{B}\left(x_{i+1}, r_{i+1}\right) \subset B\left(x_{1}, r_{1}\right) \backslash \bigcup_{j=1}^{i} \overline{E_{j}} .
$$

In particular, for each $i \in \mathbb{N}$

$$
x \notin \bigcup_{j=1}^{i} \overline{E_{j}}
$$

and therefore

$$
x \notin \bigcup_{j=1}^{\infty} \overline{E_{j}}=X
$$

a contradiction.
3.34. Definition. A function $f: X \rightarrow Y$ where $(X, \rho)$ and $(Y, \sigma)$ are metric spaces is said to bounded if there exists $0<M<\infty$ such that $\sigma(f(x), f(y)) \leq M$ for all $x, y \in X$. A family $\mathcal{F}$ of functions $f: X \rightarrow Y$ is called uniformly bounded if $\sigma(f(x), f(y)) \leq M$ for all $x, y \in X$ and for all $f \in \mathcal{F}$.

An immediate consequence of the Baire Category Theorem is the following result, which is known as the uniform boundedness principle. We will encounter this result again in the framework of functional analysis, Theorem 8.21. It states that if the upper envelope of a family of continuous functions on a complete metric space is finite everywhere, then the upper envelope is bounded above by some constant on some nonempty open subset. In other words, the family is uniformly bounded on some open set. Of course, there is no estimate of how large the open set is, but in some applications just the knowledge that such an open set exists, no matter how small, is of great importance.
3.35. Theorem. Let $\mathcal{F}$ be a family of real-valued continuous functions defined on a complete metric space $X$ and suppose

$$
\begin{equation*}
f^{*}(x):=\sup _{f \in \mathcal{F}}|f(x)|<\infty \tag{3.2}
\end{equation*}
$$

for each $x \in X$. That is, for each $x \in X$, there is a constant $M_{x}$ such that

$$
f(x) \leq M_{x} \quad \text { for all } f \in \mathcal{F}
$$

Then there exist a nonempty open set $U \subset X$ and a constant $M$ such that $|f(x)| \leq$ $M$ for all $x \in U$ and all $f \in \mathcal{F}$.
3.36. Remark. Condition (3.2) states that the family $\mathcal{F}$ is bounded at each point $x \in X$; that is, the family is pointwise bounded by $M_{x}$. In applications, a difficulty arises from the possibility that $\sup _{x \in X} M_{x}=\infty$. The main thrust of the theorem is that there exist $M>0$ and an open set $U$ such that $\sup _{x \in U} M_{x} \leq M$.

Proof. For each positive integer $i$, let

$$
E_{i, f}=\{x:|f(x)| \leq i\}, E_{i}=\bigcap_{f \in \mathcal{F}} E_{i, f} .
$$

Note that $E_{i, f}$ is closed and therefore so is $E_{i}$ since $f$ is continuous. From the hypothesis, it follows that

$$
X=\bigcup_{i=1}^{\infty} E_{i} .
$$

Since $X$ is a complete metric space, the Baire Category Theorem implies that there is some set, say $E_{M}$, that is not nowhere dense. Because $E_{M}$ is closed, it must contain an open set $U$. Now for each $x \in U$, we have $|f(x)| \leq M$ for all $f \in \mathcal{F}$, which is the desired conclusion.
3.37. Example. Here is a simple example which illustrates this result. Define a sequence of functions $f_{k}:[0,1] \rightarrow \mathbb{R}$ by

$$
f_{k}(x)= \begin{cases}k^{2} x, & 0 \leq x \leq 1 / k \\ -k^{2} x+2 k, & 1 / k \leq x \leq 2 / k \\ 0, & 2 / k \leq x \leq 1\end{cases}
$$

Thus, $f_{k}(x) \leq k$ on $[0,1]$ and $f^{*}(x) \leq k$ on $[1 / k, 1]$ and so $f^{*}(x)<\infty$ for all $0 \leq x \leq 1$. The sequence $\left\{f_{k}\right\}$ is not uniformly bounded on $[0,1]$, but it is uniformly bounded on some open set $U \subset[0,1]$. Indeed, in this example, the open set $U$ can be taken as any interval $(a, b)$ where $0<a<b<1$ because the sequence $\left\{f_{k}\right\}$ is bounded by $1 / k$ on $(2 / k, 1)$.

Exercises for Section 3.4

1. Prove that a set $E$ in a metric space is nowhere dense if and only if for each open set $U$, there is an non empty open set $V \subset U$ such that $V \cap E=\emptyset$.
2. If $(X, \rho)$ is a metric space, prove that there exists a complete metric space $\left(X^{*}, \rho^{*}\right)$ in which $X$ is isometrically embedded as a dense subset.
3. Prove that the boundary of an open set (or closed set) is nowhere dense in a topological space.

### 3.5. Compactness in Metric Spaces

In topology there are various notions related to compactness including sequential compactness and the Bolzano-Weierstrass Property. The main objective of this section is to show that these concepts are equivalent in a metric space.

The concept of completeness in a metric space is very useful, but is limited to only those sequences that are Cauchy. A stronger notion called sequential compactness allows consideration of sequences that are not Cauchy. This notion is more general in the sense that it is topological, whereas completeness is meaningful only in the setting of a metric space.

There is an abundant supply of sets that are not compact. For example, the set $A:=(0,1]$ in $\mathbb{R}$ is not compact since the collection of open intervals of the form $(1 / i, 2], i=1,2, \ldots$, provides an open cover of $A$ that admits no finite subcover. On the other hand, while it is true that $[0,1]$ is compact, the proof is not obvious. The reason for this is that the definition of compactness usually is not easy to employ directly. It is best to first determine how it intertwines with other related concepts.
3.38. Definition. Definition If $(X, \rho)$ is a metric space, a set $A \subset X$ is called totally bounded if, for every $\varepsilon>0, A$ can be covered by finitely many balls of radius $\varepsilon$. A set $A$ is bounded if there is a positive number $M$ such that $\rho(x, y) \leq M$ for all $x, y \in A$. While it is true that a totally bounded set is bounded (Exercise 3.1 ), the converse is easily seen to be false; consider (iii) of Example 3.25.
3.39. Definition. A set $A \subset X$ is said to be sequentially compact if every sequence in $A$ has a subsequence that converges to a point in $A$. Also, $A$ is said to have the Bolzano-Weierstrass property if every infinite subset of $A$ has a limit point that belongs to $A$.
3.40. Theorem. If $A$ is a subset of a metric space $(X, \rho)$, the following are equivalent:
(i) $A$ is compact.
(ii) $A$ is sequentially compact.
(iii) $A$ is complete and totally bounded.
(iv) A has the Bolzano-Weierstrass property.

Proof. Beginning with (i), we shall prove that each statement implies its successor.
(i) implies (ii): Let $\left\{x_{i}\right\}$ be a sequence in $A$; that is, there is a function $f$ defined on the positive integers such that $f(i)=x_{i}$ for $i=1,2, \ldots$. Let $E$ denote the range of $f$. If $E$ has only finitely many elements, then some member of the sequence must be repeated an infinite number of times thus showing that the sequence has a convergent subsequence.

Assuming now that $E$ is infinite we proceed by contradiction and thus suppose that $\left\{x_{i}\right\}$ had no convergent subsequence. Then each element of $E$ would be isolated. That is, with each $x \in E$ there exists $r=r_{x}>0$ such that $B\left(x, r_{x}\right) \cap E=\{x\}$. This would imply that $E$ has no limit points; thus, Theorem 3.6 (viii) and (ix), (p.39), would imply that $E$ is closed and therefore compact by Proposition 3.12. However, this would lead to a contradiction since the family $\left\{B\left(x, r_{x}\right): x \in E\right\}$ is an open cover of $E$ that possesses no finite subcover; this is impossible since $E$ consists of infinitely many points.
(ii) implies (iii): The denial of (iii) leads to two possibilities: Either $A$ is not complete or it is not totally bounded. If $A$ were not complete, there would exist a fundamental sequence $\left\{x_{i}\right\}$ in $A$ that does not converge to any point in $A$. Hence, no subsequence converges for otherwise the whole sequence would converge, thus contradicting the sequential compactness of $A$.

On the other hand, suppose $A$ is not totally bounded; then there exists $\varepsilon>0$ such that $A$ cannot be covered by finitely many balls of radius $\varepsilon$. In particular, we conclude that $A$ has infinitely many elements. Now inductively choose a sequence $\left\{x_{i}\right\}$ in $A$ as follows: select $x_{1} \in A$. Then, since $A \backslash B\left(x_{1}, \varepsilon\right) \neq \emptyset$ we can choose $x_{2} \in A \backslash B\left(x_{1}, \varepsilon\right)$. Similarly, $A \backslash\left[B\left(x_{1}, \varepsilon\right) \cup B\left(x_{2}, \varepsilon\right)\right] \neq \emptyset$ and $\rho\left(x_{1}, x_{2}\right) \geq \varepsilon$. Assuming that $x_{1}, x_{2}, \ldots, x_{i-1}$ have been chosen so that $\rho\left(x_{k}, x_{j}\right) \geq \varepsilon$ when $1 \leq k<j \leq i-1$, select

$$
x_{i} \in A \backslash \bigcup_{j=1}^{i-1} B\left(x_{j}, \varepsilon\right),
$$

thus producing a sequence $\left\{x_{i}\right\}$ with $\rho\left(x_{i}, x_{j}\right) \geq \varepsilon$ whenever $i \neq j$. Clearly, $\left\{x_{i}\right\}$ has no convergent subsequence.
(iii) implies (iv): We may as well assume that $A$ has an infinite number of elements. Under the assumptions of (iii), $A$ can be covered by finite number of balls of radius 1 and therefore, at least one of them, call it $B_{1}$, contains infinitely many points of $A$. Let $x_{1}$ be one of these points. By a similar argument, there is a ball $B_{2}$ of radius $1 / 2$ such that $A \cap B_{1} \cap B_{2}$ has infinitely many elements, and thus
it contains an element $x_{2} \neq x_{1}$. Continuing this way, we find a sequence of balls $\left\{B_{i}\right\}$ with $B_{i}$ of radius $1 / i$ and mutually distinct points $x_{i}$ such that

$$
\begin{equation*}
\bigcap_{i=1}^{k} A \bigcap B_{i} \tag{3.3}
\end{equation*}
$$

is infinite for each $k=1,2, \ldots$ and therefore contains a point $x_{k}$ distinct from $\left\{x_{1}, x_{2}, \ldots, x_{k-1}\right\}$. Observe that $0<\rho\left(x_{k}, x_{l}\right)<2 / k$ whenever $l \geq k$, thus implying that $\left\{x_{k}\right\}$ is a Cauchy sequence which, by assumption, converges to some $x_{0} \in A$. It is easy to verify that $x_{0}$ is a limit point of $A$.
(iv) implies (i): Let $\left\{U_{\alpha}\right\}$ be an arbitrary open cover of $A$. First, we claim there exists $\lambda>0$ and a countable number of balls, call them $B_{1}, B_{2}, \ldots$, such that each has radius $\lambda, A$ is contained in their union and that each $B_{k}$ is contained in some $U_{\alpha}$. To establish our claim, suppose for each positive integer $i$, there is a ball, $B_{i}$, of radius $1 / i$ such that

$$
B_{i} \cap A \neq \emptyset
$$

$$
\begin{equation*}
B_{i} \text { is not contained in any } U_{\alpha} \text {. } \tag{3.4}
\end{equation*}
$$

For each positive integer $i$, select $x_{i} \in B_{i} \cap A$. Since $A$ satisfies the BolzanoWeierstrass property, the sequence $\left\{x_{i}\right\}$ possesses a limit point and therefore it has a subsequence $\left\{x_{i_{j}}\right\}$ that converges to some $x \in A$. Now $x \in U_{\alpha}$ for some $\alpha$. Since $U_{\alpha}$ is open, there exists $\varepsilon>0$ such that $B(x, \varepsilon) \subset U_{\alpha}$. If $i_{j}$ is chosen so large that $\rho\left(x_{i_{j}}, x\right)<\frac{\varepsilon}{2}$ and $\frac{1}{i_{j}}<\frac{\varepsilon}{4}$, then for $y \in B_{i_{j}}$ we have

$$
\rho(y, x) \leq \rho\left(y, x_{i_{j}}\right)+\rho\left(x_{i_{j}}, x\right)<2 \frac{\varepsilon}{4}+\frac{\varepsilon}{2}=\varepsilon
$$

which shows that $B_{i_{j}} \subset B(x, \varepsilon) \subset U_{\alpha}$, contradicting (3.4). Thus, our claim is established.

In view of our claim, $A$ can be covered by family $\mathcal{F}$ of balls of radius $\lambda$ such that each ball belongs to some $U_{\alpha}$. A finite number of these balls also covers $A$, for if not, we could proceed exactly as in the proof above of (ii) implies (iii) to construct a sequence of points $\left\{x_{i}\right\}$ in $A$ with $\rho\left(x_{i}, x_{j}\right) \geq \lambda$ whenever $i \neq j$. This leads to a contradiction since the Bolzano-Wierstrass condition on $A$ implies that $\left\{x_{i}\right\}$ possesses a limit point $x_{0} \in A$. Thus, a finite number of balls covers $A$, say $B_{1}, \ldots B_{k}$. Each $B_{i}$ is contained in some $U_{\alpha}$, say $U_{a_{i}}$ and therefore we have

$$
A \subset \bigcup_{i=1}^{k} B_{i} \subset \bigcup_{i=1}^{k} U_{\alpha_{i}}
$$

which proves that a finite number of the $U_{\alpha}$ covers $A$.
3.41. Corollary. $A$ set $A \subset \mathbb{R}^{n}$ is compact if and only if $A$ is closed and bounded.

Proof. Clearly, $A$ is bounded if it is compact. Proposition 3.13 shows that it is also closed.

Conversely, if $A$ is closed, it is complete (see Exercise 6, Section 3.3); it thus suffices to show that any bounded subset of $\mathbb{R}^{n}$ is totally bounded. (Recall that bounded sets in an arbitrary metric space are not generally totally bounded; see Exercise 1, Section 3.5.) Since any bounded set is contained in some cube

$$
Q=[-a, a]^{n}=\left\{x \in \mathbb{R}^{n}: \max \left(\left|x_{1}\right|, \ldots,\left|x_{n}\right| \leq a\right)\right\}
$$

it is sufficient to show that $Q$ is totally bounded. For this purpose, choose $\varepsilon>0$ and let $k$ be an integer such that $k>\sqrt{n} a / \varepsilon$. Then $Q$ can be expressed as the union of $k^{n}$ congruent subcubes by dividing the interval $[-a, a]$ into $k$ equal pieces. The side length of each of these subcubes is $2 a / k$ and hence the diameter of each cube is $2 \sqrt{n} a / k<2 \varepsilon$. Therefore, each cube is contained in a ball of radius $\varepsilon$ about its center.

## Exercises for Section 3.5

1. Prove that a totally bounded set in a metric space is bounded.
2. Prove that a subset $E$ of a metric space is totally bounded if and only if $\bar{E}$ is totally bounded.
3. Prove that a totally bounded metric space is separable.
4. The proof that (iv) implies (i) in Theorem 3.40 utilizes a result that needs to be emphasized. Prove: For each open cover $\mathcal{F}$ of a compact set in a metric space, there is a number $\eta>0$ with the property that if $x, y$ are any two points in $X$ with $\rho(x, y)<\eta$, then there is an open set $V \in \mathcal{F}$ such that both $x, y$ belong to $V$. The number $\eta$ is called a Lebesgue number for the covering $\mathcal{F}$.
5 . Let $\varrho: \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}$ be defined by

$$
\varrho(x, y)=\min \{|x-y|, 1\} \quad \text { for } \quad(x, y) \in \mathbb{R} \times \mathbb{R}
$$

Prove that $\varrho$ is a metric on $\mathbb{R}$. Show that closed, bounded subsets of $(\mathbb{R}, \varrho)$ need not be compact. Hint: This metric is topologically equivalent to the Euclidean metric.

### 3.6. Compactness of Product Spaces

In this section we prove Tychonoff's Theorem which states that the product of an arbitrary number of compact topological spaces is compact. This is one of the most important theorems in general topology, in particular for its applications to functional analysis.

Let $\left\{X_{\alpha}: \alpha \in A\right\}$ be a family of topological spaces and set $X=\prod_{\alpha \in A} X_{\alpha}$. Let $P_{\alpha}: X \rightarrow X_{\alpha}$ denote the projection of $X$ onto $X_{\alpha}$ for each $\alpha$. Recall that the
family of subsets of $X$ of the form $P_{\alpha}^{-1}(U)$ where $U$ is an open subset of $X_{\alpha}$ and $\alpha \in A$ is a subbase for the product topology on $X$.

The proof of Tychonoff's theorem will utilize the finite intersection property introduced in Definition 3.15 and Lemma 3.16.

In the following proof, we utilize the Hausdorff Maximal Principle, see p. 8.
3.42. Lemma. Let $\mathcal{A}$ be a family of subsets of a set $Y$ having the finite intersection property and suppose $\mathcal{A}$ is maximal with respect to the finite intersection property, i.e., no family of subsets of $Y$ that properly contains $\mathcal{A}$ has the finite intersection property. Then
(i) $\mathcal{A}$ contains all finite intersections of members of $\mathcal{A}$.
(ii) If $S \subset Y$ and $S \cap A \neq \emptyset$ for each $A \in \mathcal{A}$, then $S \in \mathcal{A}$.

Proof. To prove (i) let $\mathcal{B}$ denote the family of all finite intersections of members of $\mathcal{A}$. Then $\mathcal{A} \subset \mathcal{B}$ and $\mathcal{B}$ has the finite intersection property. Thus by the maximality of $\mathcal{A}$, it is clear that $\mathcal{A}=\mathcal{B}$.

To prove (ii), suppose $S \cap A \neq \emptyset$ for each $A \in \mathcal{A}$. Set $\mathcal{C}=\mathcal{A} \cup\{S\}$. Then, since $\mathcal{C}$ has the finite intersection property, the maximality of $\mathcal{A}$ implies that $\mathcal{C}=\mathcal{A}$.

We can now prove Tychonoff's theorem.
3.43. Theorem (Tychonoff's Product Theorem). If $\left\{X_{\alpha}: \alpha \in A\right\}$ is a family of compact topological spaces and $X=\prod_{\alpha \in A} X_{\alpha}$ with the product topology, then $X$ is compact.

Proof. Suppose $\mathcal{C}$ is a family of closed subsets of $X$ having the finite intersection property and let $\mathcal{E}$ denote the collection of all families of subsets of $X$ such that each family contains $\mathcal{C}$ and has the finite intersection property. Then $\mathcal{E}$ satisfies the conditions of the Hausdorff Maximal Principle, and hence there is a maximal element $\mathcal{B}$ of $\mathcal{E}$ in the sense that $\mathcal{B}$ is not a subset of any other member of $\mathcal{E}$.

For each $\alpha$ the family $\left\{P_{\alpha}(B): B \in \mathcal{B}\right\}$ of subsets of $X_{\alpha}$ has the finite intersection property. Since $X_{\alpha}$ is compact, there is a point $x_{\alpha} \in X_{\alpha}$ such that

$$
x_{\alpha} \in \bigcap_{B \in \mathcal{B}} \overline{P_{\alpha}(B)} .
$$

For any $\alpha \in A$, let $U_{\alpha}$ be an open subset of $X_{\alpha}$ containing $x_{\alpha}$. Then

$$
B \bigcap P_{\alpha}^{-1}\left(U_{\alpha}\right) \neq \emptyset
$$

for each $B \in \mathcal{B}$. In view of Lemma 3.42 (ii) we see that $P_{\alpha}^{-1}\left(U_{\alpha}\right) \in \mathcal{B}$. Thus by Lemma 3.42 (i), any finite intersection of sets of this form is a member of $\mathcal{B}$. It follows that any open subset of $X$ containing $x$ has a nonempty intersection with
each member of $\mathcal{B}$. Since $\mathcal{C} \subset \mathcal{B}$ and each member of $\mathcal{C}$ is closed, it follows that $x \in C$ for each $C \in \mathcal{C}$.

## Exercises for Section 3.6

1. The set of all sequences $\left\{x_{i}\right\}_{i=1}^{\infty}$ in $[0,1]$ can be written as $[0,1]^{\mathbb{N}}$. The Tychonoff Theorem asserts that the product topology on $[0,1]^{\mathbb{N}}$ is compact. Prove that the function $\varrho$ defined by

$$
\varrho\left(\left\{x_{i}\right\},\left\{y_{i}\right\}\right)=\sum_{i=1}^{\infty} \frac{1}{2^{i}}\left|x_{i}-y_{i}\right| \quad \text { for } \quad\left\{x_{i}\right\},\left\{y_{i}\right\} \in[0,1]^{\mathbb{N}}
$$

is a metric on $[0,1]^{\mathbb{N}}$ and that this metric induces the product topology on $[0,1]^{\mathbb{N}}$. Prove that every sequence of sequences in $[0,1]$ has a convergent subsequence in the metric space $\left([0,1]^{\mathbb{N}}, \varrho\right)$. This space is sometimes called the Hilbert Cube.

### 3.7. The Space of Continuous Functions

In this section we investigate an important metric space, $C(X)$, the space of continuous functions on a metric space $X$. It is shown that this space is complete. More importantly, necessary and sufficient conditions for the compactness of subsets of $C(X)$ are given.

Recall the discussion of continuity given in Theorems 3.10 and 3.29. Our discussion will be carried out in the context of functions $f: X \rightarrow Y$ where $(X, \rho)$ and $(Y, \sigma)$ are metric spaces. Continuity of $f$ at $x_{0}$ requires that points near $x_{0}$ are mapped into points near $f\left(x_{0}\right)$. We introduce the concept of "oscillation" to assist in making this idea precise.
3.44. Definition. If $f: X \rightarrow Y$ is an arbitrary mapping, then the oscillation of $f$ on a ball $B\left(x_{0}\right)$ is defined by

$$
\operatorname{osc}\left[f, B\left(x_{0}, r\right)\right]=\sup \left\{\sigma[f(x), f(y)]: x, y \in B\left(x_{0}, r\right)\right\}
$$

Thus, the oscillation of $f$ on a ball $B\left(x_{0}, r\right)$ is nothing more than the diameter of the set $f\left(B\left(x_{0}, r\right)\right)$ in $Y$. The diameter of an arbitrary set $E$ is defined as $\sup \{\sigma(x, y): x, y \in E\}$. It may possibly assume the value $+\infty$. Note that osc $\left[f, B\left(x_{0}, r\right)\right]$ is a nondecreasing function of $r$ for each point $x_{0}$.

We leave it to the reader to supply the proof of the following assertion.
3.45. Proposition. A function $f: X \rightarrow Y$ is continuous at $x_{0} \in X$ if and only if

$$
\lim _{r \rightarrow 0} \operatorname{osc}\left[f, B\left(x_{0}, r\right)\right]=0
$$

The concept of oscillation is useful in providing information concerning the set on which an arbitrary function is continuous.

For this we need the following definitions
3.46. Definition. A subset $E$ of a topological space is called a $G_{\delta}$ set if $E$ can be written as the countable intersection of open sets, and it is an $F_{\sigma}$ set if it can be written as the countable union of closed sets.
3.47. Theorem. Let $f: X \rightarrow Y$ be an arbitrary function. Then the set of points at which $f$ is continuous is a $G_{\delta}$ set.

Proof. For each integer $i$, let

$$
G_{i}=X \cap\left\{x: \inf _{r>0} \operatorname{osc}[f, B(x, r)]<1 / i\right\}
$$

From the Proposition above, we know that $f$ is continuous at $x$ if and only if $\lim _{r \rightarrow 0} \operatorname{osc}[f, B(x, r)]=0$. Therefore, the set of points at which $f$ is continuous is given by

$$
A=\bigcap_{i=1}^{\infty} G_{i} .
$$

To complete the proof we need only show that each $G_{i}$ is open. For this, observe that if $x \in G_{i}$, then there exists $r>0$ such that osc $[f, B(x, r)]<1 / i$. Now for each $y \in B(x, r)$, there exists $t>0$ such that $B(y, t) \subset B(x, r)$ and consequently,

$$
\operatorname{osc}[f, B(y, t)] \leq \operatorname{osc}[f, B(x, r)]<1 / i
$$

This implies that each point $y$ of $B(x, r)$ is an element of $G_{i}$. That is, $B(x, r) \subset G_{i}$ and since $x$ is an arbitrary point of $G_{i}$, it follows that $G_{i}$ is open.
3.48. Theorem. Let $f$ be an arbitrary function defined on $[0,1]$ and let $E:=$ $\{x \in[0,1]:$ fis continuous at $x\}$. Then $E$ cannot be the set of rational numbers in $[0,1]$.

Proof. It suffices to show that the rationals in $[0,1]$ do not constitute a $G_{\delta}$ set. If this were false, the irrationals in $[0,1]$ would be an $F_{\sigma}$ set and thus would be the union of a countable number of closed sets, each having an empty interior. Since the rationals are a countable union of closed sets (singletons, with no interiors), it would follow that $[0,1]$ is also of the first category, contrary to the Baire Category Theorem. Thus, the rationals cannot be a $G_{\delta}$ set.

Since continuity is such a fundamental notion, it is useful to know those properties that remain invariant under a continuous transformation. The following result shows that compactness is a continuous invariant.
3.49. Theorem. Suppose $X$ and $Y$ are topological spaces and $f: X \rightarrow Y$ is a continuous mapping. If $K \subset X$ is a compact set, then $f(K)$ is a compact subset of $Y$.

Proof. Let $\mathcal{F}$ be an open cover of $f(K)$; that is, the elements of $\mathcal{F}$ are open sets whose union contains $f(K)$. The continuity of $f$ implies that each $f^{-1}(U)$ is an open subset of $X$ for each $U \in \mathcal{F}$. Moreover, the collection $\left\{f^{-1}(U): U \in \mathcal{F}\right\}$ provides an open cover of $K$. Indeed, if $x \in K$, then $f(x) \in f(K)$, and therefore that $f(x) \in U$ for some $U \in \mathcal{F}$. This implies that $x \in f^{-1}(U)$. Since $K$ is compact, $\mathcal{F}$ possesses a finite subcover for $K$, say, $\left\{f^{-1}\left(U_{1}\right), \ldots, f^{-1}\left(U_{k}\right)\right\}$. From this it easily follows that the corresponding collection $\left\{U_{1}, \ldots, U_{k}\right\}$ is an open cover of $f(K)$, thus proving that $f(K)$ is compact.
3.50. Corollary. Assume that $X$ is a compact topological space and suppose $f: X \rightarrow \mathbb{R}$ is continuous. Then, $f$ attains its maximum and minimum on $X$; that is, there are points $x_{1}, x_{2} \in X$ such that $f\left(x_{1}\right) \leq f(x) \leq f\left(x_{2}\right)$ for all $x \in X$.

Proof. From the preceding result and Corollary 3.41, it follows that $f(X)$ is a closed and bounded subset of $\mathbb{R}$. Consequently, by Theorem $2.22, f(X)$ has a least upper bound, say $y_{0}$, that belongs to $f(X)$ since $f(X)$ is closed. Thus there is a point, $x_{2} \in X$, such that $f\left(x_{2}\right)=y_{0}$. Then $f(x) \leq f\left(x_{2}\right)$ for all $x \in X$. Similarly, there is a point $x_{1}$ at which $f$ attains a minimum.

We proceed to examine yet another implication of continuous mappings defined on compact spaces. The next definition sets the stage.
3.51. Definition. Suppose $X$ and $Y$ are metric spaces. A mapping $f: X \rightarrow Y$ is said to be uniformly continuous on $\mathbf{X}$ if for each $\varepsilon>0$ there exists $\delta>0$ such that $\sigma[f(x), f(y)]<\varepsilon$ whenever $x$ and $y$ are points in $X$ with $\rho(x, y)<\delta$. The important distinction between continuity and uniform continuity is that in the latter concept, the number $\delta$ depends only on $\varepsilon$ and not on $\varepsilon$ and $x$ as in continuity. An equivalent formulation of uniform continuity can be stated in terms of oscillation, which was defined in Definition 3.44, for each number $r>0$, let

$$
\omega_{f}(r):=\sup _{x \in X} \operatorname{osc}[f, B(x, r)]
$$

The function $\omega_{f}$ is called the modulus of continuity of $f$. It is not difficult to show that $f$ is uniformly continuous on $X$ provided

$$
\lim _{r \rightarrow 0} \omega_{f}(r)=0
$$

3.52. Theorem. Let $f: X \rightarrow Y$ be a continuous mapping. If $X$ is compact, then $f$ is uniformly continuous on $X$.

Proof. Choose $\varepsilon>0$. Then the collection

$$
\mathcal{F}=\left\{f^{-1}(B(y, \varepsilon)): y \in Y\right\}
$$

is an open cover of $X$. Let $\eta$ denote the Lebesgue number of this open cover (see Exercise 4, Section 3.5). Thus, for any $x \in X$, we have $B(x, \eta / 2)$ is contained in $f^{-1}(B(y, \varepsilon))$ for some $y \in Y$. This implies $\omega_{f}(\eta / 2) \leq \varepsilon$.
3.53. Definition. For $(X, \rho)$ a metric space, let

$$
\begin{equation*}
d(f, g):=\sup (|f(x)-g(x)|: x \in X) \tag{3.5}
\end{equation*}
$$

denote the distance between any two bounded, real valued functions defined on $X$. This metric is related to the notion of uniform convergence. Indeed, a sequence of bounded functions $\left\{f_{i}\right\}$ defined on $X$ is said to converge uniformly to a bounded function $f$ on $X$ provided that $d\left(f_{i}, f\right) \rightarrow 0$ as $i \rightarrow \infty$. We denote by

$$
C(X)
$$

the space of bounded, real-valued, continuous functions on $X$.
3.54. Theorem. The space $C(X)$ is complete.

Proof. Let $\left\{f_{i}\right\}$ be a Cauchy sequence in $C(X)$. Since

$$
\left|f_{i}(x)-f_{j}(x)\right| \leq d\left(f_{i}, f_{j}\right)
$$

for all $x \in X$, it follows for each $x \in X$ that $\left\{f_{i}(x)\right\}$ is a Cauchy sequence of real numbers. Therefore, $\left\{f_{i}(x)\right\}$ converges to a number, which depends on $x$ and is denoted by $f(x)$. In this way, we define a function $f$ on $X$. In order to complete the proof, we need to show that $f$ is an element of $C(X)$ and that the sequence $\left\{f_{i}\right\}$ converges to $f$ in the metric of (3.5). First, observe that $f$ is a bounded function on $X$, because for any $\varepsilon>0$, there exists an integer $N$ such that

$$
\left|f_{i}(x)-f_{j}(x)\right|<\varepsilon
$$

whenever $x \in X$ and $i, j \geq N$. Therefore,

$$
|f(x)| \leq\left|f_{N}(x)\right|+\varepsilon
$$

for all $x \in X$, thus showing that $f$ is bounded since $f_{N}$ is.
Next, we show that

$$
\begin{equation*}
\lim _{i \rightarrow \infty} d\left(f, f_{i}\right)=0 \tag{3.6}
\end{equation*}
$$

For this, let $\varepsilon>0$. Since $\left\{f_{i}\right\}$ is a Cauchy sequence in $C(X)$, there exists $N>0$ such that $d\left(f_{i}, f_{j}\right)<\varepsilon$ whenever $i, j \geq N$. That is, $\left|f_{i}(x)-f_{j}(x)\right|<\varepsilon$ for all $i, j \geq N$ and for all $x \in X$. Thus, for each $x \in X$,

$$
\left|f(x)-f_{i}(x)\right|=\lim _{j \rightarrow \infty}\left|f_{i}(x)-f_{j}(x)\right|<\varepsilon
$$

when $i>N$. This implies that $d\left(f, f_{i}\right)<\varepsilon$ for $i>N$, which establishes (3.6) as required.

Finally, it will be shown that $f$ is continuous on $X$. For this, let $x_{0} \in X$ and $\varepsilon>0$ be given. Let $f_{i}$ be a member of the sequence such that $d\left(f, f_{i}\right)<\varepsilon / 3$. Since $f_{i}$ is continuous at $x_{0}$, there is a $\delta>0$ such that $\left|f_{i}\left(x_{0}\right)-f_{i}(y)\right|<\varepsilon / 3$ when $\rho\left(x_{0}, y\right)<\delta$. Then, for all $y$ with $\rho\left(x_{0}, y\right)<\delta$, we have

$$
\begin{aligned}
\left|f\left(x_{0}\right)-f(y)\right| & \leq\left|f\left(x_{0}\right)-f_{i}\left(x_{0}\right)\right|+\left|f_{i}\left(x_{0}\right)-f_{i}(y)\right|+\left|f_{i}(y)-f(y)\right| \\
& <d\left(f, f_{i}\right)+\frac{\varepsilon}{3}+d\left(f_{i}, f\right)<\varepsilon
\end{aligned}
$$

This shows that $f$ is continuous at $x_{0}$ and the proof is complete.
3.55. Corollary. The uniform limit of a sequence of continuous functions is continuous.

Now that we have shown that $C(X)$ is complete, it is natural to inquire about other topological properties it may possess. We will close this section with an investigation of its compactness properties. We begin by examining the consequences of uniform convergence on a compact space.
3.56. Theorem. Let $\left\{f_{i}\right\}$ be a sequence of continuous functions defined on a compact metric space $X$ that converges uniformly to a function $f$. Then, for each $\varepsilon>0$, there exists $\delta>0$ such that $\omega_{f_{i}}(r)<\varepsilon$ for all positive integers $i$ and for $0<r<\delta$.

Proof. We know from Corollary 3.55 that $f$ is continuous, and Theorem 3.52 asserts that $f$ is uniformly continuous as well as each $f_{i}$. Thus, for each $i$, we know that

$$
\lim _{r \rightarrow 0} \omega_{f_{i}}(r)=0
$$

That is, for each $\varepsilon>0$ and for each $i$, there exists $\delta_{i}>0$ such that

$$
\begin{equation*}
\omega_{f_{i}}(r)<\varepsilon \quad \text { for } \quad r<\delta_{i} \tag{3.7}
\end{equation*}
$$

However, since $f_{i}$ converges uniformly to $f$, we claim that there exists $\delta>0$ independent of $f_{i}$ such that (3.7) holds with $\delta_{i}$ replaced by $\delta$. To see this, observe that
since $f$ is uniformly continuous, there exists $\delta^{\prime}>0$ such that $|f(y)-f(x)|<\varepsilon / 3$ whenever $x, y \in X$ and $\rho(x, y)<\delta^{\prime}$. Furthermore, there exists an integer $N$ such that $\left|f_{i}(z)-f(z)\right|<\varepsilon / 3$ for $i \geq N$ and for all $z \in X$. Therefore, by the triangle inequality, for each $i \geq N$, we have

$$
\begin{align*}
\left|f_{i}(x)-f_{i}(y)\right| & \leq\left|f_{i}(x)-f(x)\right|+|f(x)-f(y)|+\left|f(y)-f_{i}(y)\right|  \tag{3.8}\\
& <\frac{\varepsilon}{3}+\frac{\varepsilon}{3}+\frac{\varepsilon}{3}=\varepsilon
\end{align*}
$$

whenever $x, y \in X$ with $\rho(x, y)<\delta^{\prime}$. Consequently, if we let

$$
\delta=\min \left\{\delta_{1}, \ldots, \delta_{N-1}, \delta^{\prime}\right\}
$$

it follows from (3.7) and (3.8) that for each positive integer $i$,

$$
\left|f_{i}(x)-f_{i}(y)\right|<\varepsilon
$$

whenever $\rho(x, y)<\delta$, thus establishing our claim.
This argument shows that the functions, $f_{i}$, are not only uniformly continuous, but that the modulus of continuity of each function tends to 0 with $r$, uniformly with respect to $i$. We use this to formulate the next definition,
3.57. Definition. A family, $\mathcal{F}$, of functions defined on $X$ is called equicontinuous if for each $\varepsilon>0$ there exists $\delta>0$ such that for each $f \in \mathcal{F},|f(x)-f(y)|<\varepsilon$ whenever $\rho(x, y)<\delta$. Alternatively, $\mathcal{F}$ is equicontinuous if for each $f \in \mathcal{F}$, $\omega_{f}(r)<\varepsilon$ whenever $0<r<\delta$. Sometimes equicontinuous families are defined pointwise; see Exercise 3.14.

We are now in a position to give a characterization of compact subsets of $C(X)$ when $X$ is a compact metric space.
3.58. ThEOREM (Arzela-Ascoli). Suppose $(X, \rho)$ is a compact metric space. Then a set $\mathcal{F} \subset C(X)$ is compact if and only if $\mathcal{F}$ is closed, bounded, and equicontinuous.

Proof. Sufficiency: It suffices to show that $\mathcal{F}$ is sequentially compact. Thus, it suffices to show that an arbitrary sequence $\left\{f_{i}\right\}$ in $\mathcal{F}$ has a convergent subsequence. Since $X$ is compact it is totally bounded, and therefore separable. Let $D=\left\{x_{1}, x_{2}, \ldots\right\}$ denote a countable, dense subset. The boundedness of $\mathcal{F}$ implies that there is a number $M^{\prime}$ such that $d(f, g)<M^{\prime}$ for all $f, g \in \mathcal{F}$. In particular, if we fix an arbitrary element $f_{0} \in \mathcal{F}$, then $d\left(f_{0}, f_{i}\right)<M^{\prime}$ for all positive integers $i$. Since $\left|f_{0}(x)\right|<M^{\prime \prime}$ for some $M^{\prime \prime}>0$ and for all $x \in X,\left|f_{i}(x)\right|<M^{\prime}+M^{\prime \prime}$ for all $i$ and for all $x$.

Our first objective is to construct a sequence of functions, $\left\{g_{i}\right\}$, that is a subsequece of $\left\{f_{i}\right\}$ and that converges at each point of $D$. As a first step toward this end, observe that $\left\{f_{i}\left(x_{1}\right)\right\}$ is a sequence of real numbers that is contained in the compact interval $[-M, M]$, where $M:=M^{\prime}+M^{\prime \prime}$. It follows that this sequence of numbers has a convergent subsequence, denoted by $\left\{f_{1 i}\left(x_{1}\right)\right\}$. Note that the point $x_{1}$ determines a subsequence of functions that converges at $x_{1}$. For example, the subsequence of $\left\{f_{i}\right\}$ that converges at the point $x_{1}$ might be $f_{1}\left(x_{1}\right), f_{3}\left(x_{1}\right), f_{5}\left(x_{1}\right), \ldots$, in which case $f_{11}=f_{1}, f_{12}=f_{3}, f_{13}=f_{5}, \ldots$. Since the subsequence $\left\{f_{1 i}\right\}$ is a uniformly bounded sequence of functions, we proceed exactly as in the previous step with $f_{1 i}$ replacing $f_{i}$. Thus, since $\left\{f_{1 i}\left(x_{2}\right)\right\}$ is a bounded sequence of real numbers, it too has a convergent subsequence which we denote by $\left\{f_{2 i}\left(x_{2}\right)\right\}$. Similar to the first step, we see that $f_{2 i}$ is a sequence of functions that is a subsequence of $\left\{f_{1 i}\right\}$ which, in turn, is a subsequence of $f_{i}$. Continuing this process, the sequence $\left\{f_{2 i}\left(x_{3}\right)\right\}$ also has a convergent subsequence, denoted by $\left\{f_{3 i}\left(x_{3}\right)\right\}$. We proceed in this way and then set $g_{i}=f_{i i}$ so that $g_{i}$ is the $i^{t h}$ function occurring in the $i^{t h}$ subsequence. We have the following situation:

$$
\begin{array}{ccccccc}
f_{11} & f_{12} & f_{13} & \ldots & f_{1 i} & \ldots & \text { first subsequence } \\
f_{21} & f_{22} & f_{23} & \ldots & f_{2 i} & \ldots & \text { subsequence of previous subsequence } \\
f_{31} & f_{32} & f_{33} & \ldots & f_{3 i} & \ldots & \text { subsequence of previous subsequence } \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \\
f_{i 1} & f_{i 2} & f_{i 3} & \ldots & f_{i i} & \ldots & \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots & i^{\text {th }} \text { subsequence }
\end{array}
$$

Observe that the sequence of functions $\left\{g_{i}\right\}$ converges at each point of $D$. Indeed, $g_{i}$ is an element of the $j^{\text {th }}$ row for $i \geq j$. In other words, the tail end of $\left\{g_{i}\right\}$ is a subsequence of $\left\{f_{j i}\right\}$ for any $j \in \mathbb{N}$, and so it will converge as $i \rightarrow \infty$ at any point for which $\left\{f_{j i}\right\}$ converges as $i \rightarrow \infty$; i.e., for each point of $D$.

We now proceed to show that $\left\{g_{i}\right\}$ converges at each point of $X$ and that the convergence is, in fact, uniform on $X$. For this purpose, choose $\varepsilon>0$ and let $\delta>0$ be the number obtained from the definition of equicontinuity. Since $X$ is compact it is totally bounded, and therefore there is a finite number of balls of radius $\delta / 2$, say $k$ of them, whose union covers $X: X=\bigcup_{i=1}^{k} B_{i}(\delta / 2)$. Then selecting any $y_{i} \in B_{i}(\delta / 2) \cap D$ it follows that

$$
X=\bigcup_{i=1}^{k} B\left(y_{i}, \delta\right)
$$

Let $D^{\prime}:=\left\{y_{1}, y_{2}, \ldots, y_{k}\right\}$ and note $D^{\prime} \subset D$. Therefore each of the $k$ sequences

$$
\left\{g_{i}\left(y_{1}\right)\right\},\left\{g_{i}\left(y_{2}\right)\right\}, \ldots,\left\{g_{i}\left(y_{k}\right)\right\}
$$

converges, and so there is an integer $N \in \mathbb{N}$ such that if $i, j \geq N$, then

$$
\left|g_{i}\left(y_{m}\right)-g_{j}\left(y_{m}\right)\right|<\varepsilon \text { for } m=1,2, \ldots, k
$$

For each $x \in X$, there exists $y_{m} \in D^{\prime}$ such that $\left|x-y_{m}\right|<\delta$. Thus, by equicontinuity, it follows that

$$
\left|g_{i}(x)-g_{i}\left(y_{m}\right)\right|<\varepsilon
$$

for all positive integers $i$. Therefore, we have

$$
\begin{aligned}
\left|g_{i}(x)-g_{j}(x)\right| \leq & \left|g_{i}(x)-g_{i}\left(y_{m}\right)\right|+\left|g_{i}\left(y_{m}\right)-g_{j}\left(y_{m}\right)\right| \\
& +\left|g_{j}\left(y_{m}\right)-g_{j}(x)\right| \\
< & \varepsilon+\varepsilon+\varepsilon=3 \varepsilon
\end{aligned}
$$

provided $i, j \geq N$. This shows that

$$
d\left(g_{i}, g_{j}\right)<3 \varepsilon \quad \text { for } \quad i, j \geq N
$$

That is, $\left\{g_{i}\right\}$ is a Cauchy sequence in $\mathcal{F}$. Since $C(X)$ is complete (Theorem 3.54) and $\mathcal{F}$ is closed, it follows that $\left\{g_{i}\right\}$ converges to an element $g \in \mathcal{F}$. Since $\left\{g_{i}\right\}$ is a subsequence of the original sequence $\left\{f_{i}\right\}$, we have shown that $\mathcal{F}$ is sequentially compact, thus establishing the sufficiency argument.

Necessity: Note that $\mathcal{F}$ is closed since $\mathcal{F}$ is assumed to be compact. Furthermore, the compactness of $\mathcal{F}$ implies that $\mathcal{F}$ is totally bounded and therefore bounded. For the proof that $\mathcal{F}$ is equicontinuous, note that $\mathcal{F}$ being totally bounded implies that for each $\varepsilon>0$, there exist a finite number of elements in $\mathcal{F}$, say $f_{1}, \ldots, f_{k}$, such that any $f \in \mathcal{F}$ is within $\varepsilon / 3$ of $f_{i}$, for some $i \in\{1, \ldots, k\}$. Consequently, by Exercise 3.5, we have

$$
\begin{equation*}
\omega_{f}(r) \leq \omega_{f_{i}}(r)+2 d\left(f, f_{i}\right)<\omega_{f_{i}}(r)+2 \varepsilon / 3 . \tag{3.9}
\end{equation*}
$$

Since $X$ is compact, each $f_{i}$ is uniformly continuous on $X$. Thus, for each $i, i=$ $1, \ldots, k$, there exists $\delta_{i}>0$ such that $\omega_{f_{i}}(r)<\varepsilon / 3$ for $r<\delta_{i}$. Now let $\delta=$ $\min \left\{\delta_{1}, \ldots, \delta_{k}\right\}$. By (3.9) it follows that $\omega_{f}(r)<\varepsilon$ whenever $r<\delta$, which proves that $\mathcal{F}$ is equicontinuous.

In many applications, it is not of great interest to know whether $\mathcal{F}$ itself is compact, but whether a given sequence in $\mathcal{F}$ has a subsequence that converges uniformly to an element of $C(X)$, and not necessarily to an element of $\mathcal{F}$. In other words, the compactness of the closure of $\mathcal{F}$ is the critical question. It is easy to see that if $\mathcal{F}$ is equicontinuous, then so is $\overline{\mathcal{F}}$. This leads to the following corollary.
3.59. Corollary. Suppose $(X, \rho)$ is a compact metric space and suppose that $\mathcal{F} \subset C(X)$ is bounded and equicontinuous. Then $\overline{\mathcal{F}}$ is compact.

Proof. This follows immediately from the previous theorem since $\overline{\mathcal{F}}$ is both bounded and equicontinuous.

In particular, this corollary yields the following special result.
3.60. Corollary. Let $\left\{f_{i}\right\}$ be an equicontinuous, uniformly bounded sequence of functions defined on $[0,1]$. Then there is a subsequence that converges uniformly to a continuous function on $[0,1]$.

We close this section with a result that will be used frequently throughout the sequel.
3.61. Theorem. Suppose $f$ is a bounded function on $[a, b]$ that is either nondecreasing or nonincreasing. Then $f$ has at most a countable number of discontinuities.

Proof. We will give the proof only in case $f$ is nondecreasing, the proof for $f$ nonincreasing being essentially the same.

Since $f$ is nondecreasing, it follows that the left and right-hand limits exist at each point (see Exercise 25, Section 3.7) and the discontinuities of $f$ occur precisely where these limits are not equal. Thus, setting

$$
f\left(x^{+}\right)=\lim _{y \rightarrow x^{+}} f(y) \quad \text { and } \quad f\left(x^{-}\right)=\lim _{y \rightarrow x^{-}} f(y)
$$

the set $D$ of discontinuities of $f$ in $(a, b)$ is given by

$$
D=(a, b) \cap\left(\bigcup_{k=1}^{\infty}\left\{x: f\left(x^{+}\right)-f\left(x^{-}\right)>\frac{1}{k}\right\}\right) .
$$

For each $k$ the set

$$
\left\{x: f\left(x^{+}\right)-f\left(x^{-}\right)>\frac{1}{k}\right\}
$$

is finite since $f$ is bounded and thus $D$ is countable.

## Exercises for Section 3.7

1. Prove that the set of rational numbers in the real line is not a $G_{\delta}$ set.
2. Prove that the two definitions of uniform continuity given in Definition 3.51 are equivalent.
3. Assume that $(X, \rho)$ is a metric space with the property that each function $f: X \rightarrow \mathbb{R}$ is uniformly continuous.
(a) Show that $X$ is a complete metric space
(b) Give an example of a space $X$ with the above property that is not compact.
(c) Prove that if $X$ has only a finite number of isolated points, then $X$ is compact. See p. 54 for the definition of isolated point.
4. Prove that a family of functions $F$ is equicontinuous provided there exists a nondecreasing real valued function $\varphi$ such that

$$
\lim _{r \rightarrow 0} \varphi(r)=0
$$

and $\omega_{f}(r) \leq \varphi(r)$ for all $f \in F$.
5. Suppose $f, g$ are any two functions defined on a metric space. Prove that

$$
\omega_{f}(r) \leq \omega_{g}(r)+2 d(f, g)
$$

6. Prove that a Lipschitzian function is uniformly continuous.
7. Prove: If $F$ is a family of Lipschitzian functions from a bounded metric space $X$ into a metric space $Y$ such that $M$ is a Lipschitz constant for each member of $F$ and $\left\{f\left(x_{0}\right): f \in F\right\}$ is a bounded set in $Y$ for some $x_{0} \in X$, then $F$ is a uniformly bounded, equicontinuous family.
8. Let $(X, \varrho)$ and $(Y, \sigma)$ be metric spaces and let $f: X \rightarrow Y$ be uniformly continuous. Prove that if $X$ is totally bounded, then $f(X)$ is totally bounded.
9. Let $(X, \varrho)$ and $(Y, \sigma)$ be metric spaces and let $f: X \rightarrow Y$ be an arbitrary function. The graph of $f$ is a subset of $X \times Y$ defined by

$$
G_{f}:=\{(x, y): y=f(x)\} .
$$

Let $d$ be the metric $d_{1}$ on $X \times Y$ as defined in Exercise 4, Section 3.3. If $Y$ is compact, show that $f$ is continuous if and only if $G_{f}$ is a closed subset of the metric space $(X \times Y, d)$. Can the compactness assumption on $Y$ be dropped?
10. Let $Y$ be a dense subset of a metric space $(X, \varrho)$. Let $f: Y \rightarrow Z$ be a uniformly continuous function where $Z$ is a complete metric space. Show that there is a uniformly continuous function $g: X \rightarrow Z$ with the property that $f=g\llcorner Y$. Can the assumption of uniform continuity be relaxed to mere continuity?
11. Exhibit a bounded function that is continuous on $(0,1)$ but not uniformly continuous.
12. Let $\left\{f_{i}\right\}$ be a sequence of real-valued, uniformly continuous functions on a metric space $(X, \rho)$ with the property that for some $M>0,\left|f_{i}(x)-f_{j}(x)\right| \leq M$ for all positive integers $i, j$ and all $x \in X$. Suppose also that $d\left(f_{i}, f_{j}\right) \rightarrow 0$ as $i, j \rightarrow \infty$. Prove that there is a uniformly continuous function $f$ on $X$ such that $d\left(f_{i}, f\right) \rightarrow 0$ as $i \rightarrow \infty$.
13. Let $X \xrightarrow{f} Y$ where $(X, \rho)$ and $(Y, \sigma)$ are metric spaces and where $f$ is continuous. Suppose $f$ has the following property: For each $\varepsilon>0$ there is a compact set
$K_{\varepsilon} \subset X$ such that $\sigma(f(x), f(y))<\varepsilon$ for all $x, y \in X \backslash K_{e}$. Prove that $f$ is uniformly continuous on $X$.
14. A family $\mathcal{F}$ of functions defined on a metric space $X$ is called equicontinuous at $x \in X$ if for every $\varepsilon>0$ there exists $\delta>0$ such that $|f(x)-f(y)|<\varepsilon$ for all $y$ with $|x-y|<\delta$ and all $f \in \mathcal{F}$. Show that the Arzela-Ascoli Theorem remains valid with this definition of equicontinuity. That is, prove that if $X$ is compact and $\mathcal{F}$ is closed, bounded, and equicontinuous at each $x \in X$, then $\mathcal{F}$ is compact.
15. Give an example of a sequence of real valued functions defined on $[a, b]$ that converges uniformly to a continuous function, but is not equicontinuous.
16. Let $\left\{f_{i}\right\}$ be a sequence of nonnegative, equicontinuous functions defined on a totally bounded metric space $X$ such that

$$
\limsup _{i \rightarrow \infty} f_{i}(x)<\infty \text { for each } x \in X
$$

Prove that there is a subsequence that converges uniformly to a continuous function $f$.
17. Let $\left\{f_{i}\right\}$ be a sequence of nonnegative, equicontinuous functions defined on $[0,1]$ with the property that

$$
\limsup _{i \rightarrow \infty} f_{i}\left(x_{0}\right)<\infty \text { for some } x_{0} \in[0,1]
$$

Prove that there is a subsequence that converges uniformly to a continuous function $f$.
18. Let $\left\{f_{i}\right\}$ be a sequence of nonnegative, equicontinuous functions defined on a locallly compact metric space $X$ such that

$$
\limsup _{i \rightarrow \infty} f_{i}(x)<\infty \text { for each } x \in X
$$

Prove that there is an open set $U$ and a subsequence that converges uniformly on $U$ to a continuous function $f$.
19. Let $\left\{f_{i}\right\}$ be a sequence of real valued functions defined on a compact metric space $X$ with the property that $x_{k} \rightarrow x$ implies $f_{k}\left(x_{k}\right) \rightarrow f(x)$ where $f$ is a continuous function on $X$. Prove that $f_{k} \rightarrow f$ uniformly on $X$.
20. Let $\left\{f_{i}\right\}$ be a sequence of nondecreasing, real valued (not necessarily continuous) functions defined on $[a, b]$ that converges pointwise to a continuous function $f$. Show that the convergence is necessarily uniform.
21. Let $\left\{f_{i}\right\}$ be a sequence of continuous, real valued functions defined on a compact metric space $X$ that converges pointwise on some dense set to a continuous function on $X$. Prove that $f_{i} \rightarrow f$ uniformly on $X$.
22. Let $\left\{f_{i}\right\}$ be a uniformly bounded sequence in $C[a, b]$. For each $x \in[a, b]$, define

$$
F_{i}(x):=\int_{a}^{x} f_{i}(t) d t
$$

Prove that there is a subsequence of $\left\{F_{i}\right\}$ that converges uniformly to some function $F \in C[a, b]$.
23. For each integer $k>1$ let $\mathcal{F}_{k}$ be the family of continuous functions on $[0,1]$ with the property that for some $x \in[0,1-1 / k]$ we have

$$
|f(x+h)-f(x)| \leq k h \quad \text { whenever } 0<h<\frac{1}{k}
$$

(a) Prove that $\mathcal{F}_{k}$ is nowhere dense in the space $C[0,1]$ endowed with its usual metric of uniform convergence.
(b) Using the Baire Category theorem, prove that there exists $f \in C[0,1]$ that is not differentiable at any point of $(0,1)$.
24. The previous problem demonstrates the remarkable fact that functions that are nowhere differentiable are in great abundance whereas functions that are wellbehaved are relatively scarce. The following are examples of functions that are continuous and nowhere differentiable.
(a) For $x \in[0,1]$ let

$$
f(x):=\sum_{n=0}^{\infty} \frac{\left[10^{n} x\right]}{10^{n}}
$$

where [ $y$ ] denotes the distance from the greatest integer in $y$.
(b)

$$
f(x):=\sum_{k=0}^{\infty} a^{k} \cos ^{b} \pi x
$$

where $1<a b<b$. Weierstrass was the first to prove the existence of continuous nowhere differentiable functions by conceiving of this function and then proving that it is nowhere differentiable for certain values of $a$ and $b$, [49]. Later, Hardy proved the same result for all $a$ and $b,[34]$.
25. Let $f$ be a non-decreasing function on $(a, b)$. Show that $f(x+)$ and $f(x-)$ exist at every point of $x$ of $(a, b)$. Show also that if $a<x<y<b$ then $f(x+) \leq f(y-)$.

### 3.8. Lower Semicontinuous Functions

In many applications in analysis, lower and upper semicontinuous functions play an important role. The purpose of this section is to introduce these functions and develop their basic properties.

Recall that a function $f$ on a metric space is continuous at $x_{0}$ if for each $\varepsilon>0$, there exists $r>0$ such that

$$
f\left(x_{0}\right)-\varepsilon<f(x)<f\left(x_{0}\right)+\varepsilon
$$

whenever $x \in B\left(x_{0}, r\right)$. Semicontinuous functions require only one part of this inequality to hold.
3.62. Definition. Suppose $(X, \rho)$ is a metric space. A function $f$ defined on $X$ with possibly infinite values is said to be lower semicontinuous at $x_{0} \in X$ if the following conditions hold. If $f\left(x_{0}\right)<\infty$, then for every $\varepsilon>0$ there exists $r>0$ such that $f(x)>f\left(x_{0}\right)-\varepsilon$ whenever $x \in B\left(x_{0}, r\right)$. If $f\left(x_{0}\right)=\infty$, then for every positive number $M$ there exists $r>0$ such that $f(x) \geq M$ for all $x \in B\left(x_{0}, r\right)$. The function $f$ is called lower semicontinuous if it is lower semicontinuous at all $x \in X$. An upper semicontinuous function is defined analogously: if $f\left(x_{0}\right)>-\infty$, then $f(x)<f\left(x_{0}\right)+\varepsilon$ for all $x \in B\left(x_{0}, r\right)$. If $f\left(x_{0}\right)=-\infty$, then $f(x)<-M$ for all $x \in B\left(x_{0}, r\right)$.

Of course, a continuous function is both lower and upper semicontinuous. It is easy to see that the characteristic function of an open set is lower semicontinuous and that the characteristic function of a closed set is upper semicontinuous.

Semicontinuity can be reformulated in terms of the lower limit (also called limit inferior) and upper limit of (also called limit superior) $f$.
3.63. Definition. We define

$$
\liminf _{x \rightarrow x_{0}} f(x)=\lim _{r \rightarrow 0} m\left(r, x_{0}\right)
$$

where $m\left(r, x_{0}\right)=\inf \left\{f(x): 0<\rho\left(x, x_{0}\right)<r\right\}$. Similarly,

$$
\limsup _{x \rightarrow x_{0}} f(x)=\lim _{r \rightarrow 0} M\left(r, x_{0}\right)
$$

where $M\left(r, x_{0}\right)=\sup \left\{f(x): 0<\rho\left(x, x_{0}\right)<r\right\}$.
One readily verifies that $f$ is lower semicontinuous at a limit point $x_{0}$ of $X$ if and only if

$$
\liminf _{x \rightarrow x_{0}} f(x) \geq f\left(x_{0}\right)
$$

and $f$ is upper semicontinuous at $x_{0}$ if and only if

$$
\limsup _{x \rightarrow x_{0}} f(x) \leq f\left(x_{0}\right)
$$

In terms of sequences, these statements are equivalent, respectively, to the following:

$$
\liminf _{k \rightarrow \infty} f\left(x_{k}\right) \geq f\left(x_{0}\right)
$$

and

$$
\limsup _{k \rightarrow \infty} f\left(x_{k}\right) \leq f\left(x_{0}\right)
$$

whenever $\left\{x_{k}\right\}$ is a sequence converging to $x_{0}$. This leads immediately to the following.
3.64. Theorem. Suppose $X$ is a compact metric space. Then a real valued lower (upper) semicontinuous function on $X$ assumes its minimum (maximum) on $X$.

Proof. We will give the proof for $f$ lower semicontinuous, the proof for $f$ upper semicontinuous being similar. Let

$$
m=\inf \{f(x): x \in X\}
$$

We will see that $m \neq-\infty$ and that there exists $x_{0} \in X$ such that $f\left(x_{0}\right)=m$, thus establishing the result.

To see this, let $y_{k} \in f(X)$ such that $\left\{y_{k}\right\} \rightarrow m$ as $k \rightarrow \infty$. At this point of the proof, we must allow the possibility that $m=-\infty$. Note that $m \neq+\infty$. Let $x_{k} \in X$ be such that $f\left(x_{k}\right)=y_{k}$. Since $X$ is compact, there is a point $x_{0} \in X$ and a subsequence (still denoted by $\left\{x_{k}\right\}$ ) such that $\left\{x_{k}\right\} \rightarrow x_{0}$. Since $f$ is lower semicontinuous, we obtain

$$
m=\liminf _{k \rightarrow \infty} f\left(x_{k}\right) \geq f\left(x_{0}\right)
$$

which implies that $f\left(x_{0}\right)=m$ and that $m \neq-\infty$.
The following result will require the definition of a Lipschitz function.
3.65. Definition. Suppose $(X, \rho)$ and $(Y, \sigma)$ are metric spaces. A mapping $f: X \rightarrow Y$ is called Lipschitz if there is a constant $C_{f}$ such that

$$
\begin{equation*}
\sigma[f(x), f(y)] \leq C_{f} \rho(x, y) \tag{3.10}
\end{equation*}
$$

for all $x, y \in X$. The smallest such constant $C_{f}$ is called the Lipschitz constant of $f$.
3.66. Theorem. Suppose $(X, \rho)$ is a metric space.
(i) $f$ is lower semicontinuous on $X$ if and only if $\{f>t\}$ is open for all $t \in \mathbb{R}$.
(ii) If both $f$ and $g$ are lower semicontinuous on $X$, then $\min \{f, g\}$ is lower semicontinuous.
(iii) The upper envelope of any collection of lower semicontinuous functions is lower semicontinuous.
(iv) Each nonnegative lower semicontinuous function on $X$ is the upper envelope of a nondecreasing sequence of continuous (in fact, Lipschitzian) functions.

Proof. To prove (i), choose $x_{0} \in\{f>t\}$. Let $\varepsilon=f\left(x_{0}\right)-t$, and use the definition of lower semicontinuity to find a ball $B\left(x_{0}, r\right)$ such that $f(x)>f\left(x_{0}\right)-\varepsilon=$ $t$ for all $x \in B\left(x_{0}, r\right)$. Thus, $B\left(x_{0}, r\right) \subset\{f>t\}$, which proves that $\{f>t\}$ is open.

Conversely, choose $x_{0} \in X$ and $\varepsilon>0$ and let $t=f\left(x_{0}\right)-\varepsilon$. Then $x_{0} \in\{f>t\}$ and since $\{f>t\}$ is open, there exists a ball $B\left(x_{0}, r\right) \subset\{f>t\}$. This implies that $f(x)>f\left(x_{0}\right)-\varepsilon$ whenever $x \in B\left(x_{0}, r\right)$, thus establishing lower semicontinuity.
(i) immediately implies (ii) and (iii). For (ii), let $h=\min (f, g)$ and observe that $\{h>t\}=\{f>t\} \cap\{g>t\}$, which is the intersection of two open sets.

Similarly, for (iii) let $\mathcal{F}$ be a family of lower semicontinuous functions and set

$$
h(x)=\sup \{f(x): f \in \mathcal{F}\} \quad \text { for } \quad x \in X
$$

Then, for each real number $t$,

$$
\{h>t\}=\bigcup_{f \in \mathcal{F}}\{f>t\}
$$

which is open since each set on the right is open.
Proof of (iv): For each positive integer $k$ define

$$
f_{k}(x)=\inf \{f(y)+k \rho(x, y): y \in X\} .
$$

Observe that $f_{1} \leq f_{2}, \ldots, \leq f$. To show that each $f_{k}$ is Lipschitzian, it is sufficient to prove

$$
\begin{equation*}
f_{k}(x) \leq f_{k}(w)+k \rho(x, w) \quad \text { for all } \quad w \in X \tag{3.11}
\end{equation*}
$$

since the roles of $x$ and $w$ can be interchanged. To prove (3.11) observe that for each $\varepsilon>0$, there exists $y \in X$ such that

$$
f_{k}(w) \leq f(y)+k \rho(w, y) \leq f_{k}(w)+\varepsilon
$$

Now,

$$
\begin{aligned}
f_{k}(x) & \leq f(y)+k \rho(x, y) \\
& =f(y)+k \rho(w, y)+k \rho(x, y)-k \rho(w, y) \\
& \leq f_{k}(w)+\varepsilon+k \rho(x, w)
\end{aligned}
$$

where the triangle inequality has been used to obtain the last inequality. This implies (3.11) since $\varepsilon$ is arbitrary.

Finally, to show that $f_{k}(x) \rightarrow f(x)$ for each $x \in X$, observe for each $x \in X$ there is a sequence $\left\{x_{k}\right\} \subset X$ such that

$$
f\left(x_{k}\right)+k \rho\left(x_{k}, x\right) \leq f_{k}(x)+\frac{1}{k} \leq f(x)+1<\infty
$$

As a consequence, we have that $\lim _{k \rightarrow \infty} \rho\left(x_{k}, x\right)=0$. Given $\varepsilon>0$ there exists $n \in \mathbb{N}$ such that

$$
f_{k}(x)+\varepsilon \geq f_{k}(x)+\frac{1}{k} \geq f\left(x_{k}\right) \geq f(x)-\varepsilon
$$

whenever $k \geq n$ and thus $f_{k}(x) \rightarrow f(x)$.
3.67. Remark. Of course, the previous theorem has a companion that pertains to upper semicontinuous functions. Thus, the result analogous to (i) states that $f$ is upper semicontinuous on $X$ if and only if $\{f<t\}$ is open for all $t \in \mathbb{R}$. We leave it to the reader to formulate and prove the remaining three statements.
3.68. Definition. Theorem 3.66 provides a means of defining upper and lower semicontinuity for functions defined merely on a topological space $X$.

Thus, $f: X \rightarrow \mathbb{R}$ is called upper semicontinuous (lower semicontinuous) if $\{f<t\}(\{f>t\})$ is open for all $t \in \mathbb{R}$. It is easily verified that (ii) and (iii) of Theorem 3.66 remain true when $X$ is assumed to be only a topological space.

## Exercises for Section 3.8

1. Let $\left\{f_{i}\right\}$ be a decreasing sequence of upper semicontinuous functions defined on a compact metric space $X$ such that $f_{i}(x) \rightarrow f(x)$ where $f$ is lower semicontinuous. Prove that $f_{i} \rightarrow f$ uniformly.
2. Show that Theorem 3.66, (iv), remains true for lower semicontinuous functions that are bounded below. Show also that this assumption is necessary.

## CHAPTER 4

## Measure Theory

### 4.1. Outer Measure


#### Abstract

An outer measure on an abstract set $X$ is a monotone, countably subadditive function defined on all subsets of $X$. In this section, the notion of measurable set is introduced, and it is shown that the class of measurable sets forms a $\sigma$-algebra, i.e., measurable sets are closed under the operations of complementation and countable unions. It is also shown that an outer measure is countably additive on disjoint measurable sets.


In this section we introduce the concept of outer measure that will underlie and motivate some of the most important concepts of abstract measure theory. The "length" of set in $\mathbb{R}$, the "area" of a set in $\mathbb{R}^{2}$ or the "volume" of a set in $\mathbb{R}^{3}$ are notions that can be developed from basic and strongly intuitive geometric principles provided the sets are well-behaved. If one wished to develop a concept of volume in $\mathbb{R}^{3}$, for example, that would allow the assignment of volume to any set, then one could hope for a function $\mathcal{V}$ that assigns to each subset $E \subset \mathbb{R}^{3}$ a number $\mathcal{V}(E) \in[0, \infty]$ having the following properties:
(i) If $\left\{E_{i}\right\}_{i=1}^{k}$ is any finite sequence of mutually disjoint sets, then

$$
\begin{equation*}
\mathcal{V}\left(\bigcup_{i=1}^{k} E_{i}\right)=\sum_{i=i}^{k} \mathcal{V}\left(E_{i}\right) \tag{4.1}
\end{equation*}
$$

(ii) If two sets $E$ and $F$ are congruent, then $\mathcal{V}(E)=\mathcal{V}(F)$
(iii) $\mathcal{V}(Q)=1$ where $Q$ is the cube of side length 1 .

However, these three conditions are inconsistent. In 1924, Banach and Tarski [2] proved that it is possible to decompose a ball in $\mathbb{R}^{3}$ into six pieces which can be reassembled by rigid motions to form two balls, each the same size as the original. The sets in this decomposition are pathological and require the axiom of choice for their existence. If condition (i) is changed to require countable additivity rather than mere finite additivity; that is, require that if $\left\{E_{i}\right\}_{i=1}^{\infty}$ is any infinite sequence of mutually disjoint sets, then

$$
\sum_{i=i}^{\infty} \mathcal{V}\left(E_{i}\right)=\mathcal{V}\left(\bigcup_{i=1}^{\infty} E_{i}\right)
$$

this too suffers from the same inconsistency, and thus we are led to the conclusion that there is no function $\mathcal{V}$ satisfying all three conditions above. Later, we will also see that if we restrict $\mathcal{V}$ to a large class of subsets of $\mathbb{R}^{3}$ that omits only the truly pathological sets, then it is possible to incorporate $\mathcal{V}$ in a satisfactory theory of volume.

We will proceed to find this large class of sets by considering a very general context and replace countable additivity by countable subadditivity.
4.1. Definition. A function $\varphi$ defined for every subset $A$ of an arbitrary set $X$ is called an outer measure on $X$ provided the following conditions are satisfied:
(i) $\varphi(\emptyset)=0$,
(ii) $0 \leq \varphi(A) \leq \infty$ whenever $A \subset X$,
(iii) $\varphi\left(A_{1}\right) \leq \varphi\left(A_{2}\right)$ whenever $A_{1} \subset A_{2}$,
(iv) $\varphi\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \sum_{i=1}^{\infty} \varphi\left(A_{i}\right)$ for any countable collection of sets $\left\{A_{i}\right\}$ in $X$.

Condition (iii) states that $\varphi$ is monotone while (iv) states that $\varphi$ is countably subadditive. As we mentioned earlier, suitable additivity properties are necessary in measure theory; subadditivity, in general, will not suffice to produce a useful theory. We will now introduce the concept of a "measurable set" and show later that measurable sets enjoy a wide spectrum of additivity properties.

The term "outer measure" is derived from the way outer measures are constructed in practice. Often one uses a set function that is defined on some family of primitive sets (such as the family of intervals in $\mathbb{R}$ ) to approximate an arbitrary set from the "outside" to define its measure. Examples of this procedure will be given in Sections 4.3 and 4.4. First, consider some elementary examples of outer measures.
4.2. Examples. (i) In an arbitrary set $X$, define $\varphi(A)=1$ if $A$ is nonempty and $\varphi(\emptyset)=0$.
(ii) Let $\varphi(A)$ be the number (possibly infinite) of points in $A$.
(iii) Let $\varphi(\mathrm{A})= \begin{cases}0 & \text { if } \operatorname{card} A \leq \aleph_{0} \\ 1 & \text { if } \operatorname{card} A>\aleph_{0}\end{cases}$
(iv) If $X$ is a metric space, fix $\varepsilon>0$. Let $\varphi(A)$ be the smallest number of balls of radius $\varepsilon$ that cover $A$.
(v) Select a fixed $x_{0}$ in an arbitrary set $X$, and let

$$
\varphi(A)= \begin{cases}0 & \text { if } x_{0} \notin A \\ 1 & \text { if } x_{0} \in A\end{cases}
$$

$\varphi$ is called the Dirac measure concentrated at $x_{0}$.

Notice that the domain of an outer measure $\varphi$ is $\mathcal{P}(X)$, the collection of all subsets of $X$. In general it may happen that the equality $\varphi(A \cup B)=\varphi(A)+\varphi(B)$ fails when $A \cap B=\emptyset$. This property and more generally, property (4.1), will require a more restrictive class of subsets of $X$, called measurable sets, which we now define.
4.3. Definition. Let $\varphi$ be an outer measure on a set $X$. A set $E \subset X$ is called $\varphi$-measurable if

$$
\varphi(A)=\varphi(A \cap E)+\varphi(A-E)
$$

for every set $A \subset X$. In view of property (iv) above, observe that $\varphi$-measurability only requires

$$
\begin{equation*}
\varphi(A) \geq \varphi(A \cap E)+\varphi(A-E) \tag{4.2}
\end{equation*}
$$

This definition, while not very intuitive, says that a set is $\varphi$-measurable if it decomposes an arbitrary set, $A$, into two parts for which $\varphi$ is additive. We use this definition in deference to Carathéodory, who established this property as an alternative characterization of measurability in the special case of Lebesgue measure (see Definition 4.21 below). The following characterization of $\varphi$-measurability is perhaps more intuitively appealing.
4.4. Lemma. $A$ set $E \subset X$ is $\varphi$-measurable if and only if

$$
\varphi(P \cup Q)=\varphi(P)+\varphi(Q)
$$

for any sets $P$ and $Q$ such that $P \subset E$ and $Q \subset \tilde{E}$.
Proof. Sufficiency: Let $A \subset X$. Then with $P:=A \cap E \subset E$ and $Q:=$ $A-E \subset \tilde{E}$ we have $A=P \cup Q$ and therefore

$$
\varphi(A)=\varphi(P \cup Q)=\varphi(P)+\varphi(Q)=\varphi(A \cap E)+\varphi(A-E)
$$

Necessity: Let $P$ and $Q$ be arbitrary sets such that $P \subset E$ and $Q \subset \tilde{E}$. Then, by the definition of $\varphi$-measurability,

$$
\begin{aligned}
\varphi(P \cup Q) & =\varphi[(P \cup Q) \cap E]+\varphi[(P \cup Q) \cap \tilde{E}] \\
& =\varphi(P \cap E)+\varphi(Q \cap \tilde{E}) \\
& =\varphi(P)+\varphi(Q) .
\end{aligned}
$$

4.5. Remark. Recalling Examples 4.2, one verifies that only the empty set and $X$ are measurable for (i) while all sets are measurable for (ii).

Now that we have an alternate definition of $\varphi$-measurability, we investigate the properties of $\varphi$-measurable sets. We start with the following theorem which is basic
to the theory. A set function that satisfies property (iv) below on any sequence of disjoint sets is said to be countably additive.
4.6. Theorem. Suppose $\varphi$ is an outer measure on an arbitrary set $X$. Then the following four statements hold.
(i) $E$ is $\varphi$-measurable whenever $\varphi(E)=0$,
(ii) $\emptyset$ and $X$ are $\varphi$-measurable.
(iii) $E_{1}-E_{2}$ is $\varphi$-measurable whenever $E_{1}$ and $E_{2}$ are $\varphi$-measurable,
(iv) If $\left\{E_{i}\right\}$ is a countable collection of disjoint $\varphi$-measurable sets, then $\cup_{i=1}^{\infty} E_{i}$ is $\varphi$-measurable and

$$
\varphi\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\sum_{i=1}^{\infty} \varphi\left(E_{i}\right)
$$

More generally, if $A \subset X$ is an arbitrary set, then

$$
\varphi(A)=\sum_{i=1}^{\infty} \varphi\left(A \cap E_{i}\right)+\varphi(A \cap \tilde{S})
$$

where $S=\bigcup_{i=1}^{\infty} E_{i}$.
Proof. (i) If $A \subset X$, then $\varphi(A \cap E)=0$. Thus, $\varphi(A) \leq \varphi(A \cap E)+\varphi(A \cap$ $\tilde{E})=\varphi(A \cap \tilde{E}) \leq \varphi(A)$.
(ii) This follows immediately from Lemma 4.4.
(iii) We will use Lemma 4.4 to establish the $\varphi$-measurability of $E_{1}-E_{2}$. Thus, let $P \subset E_{1}-E_{2}$ and $Q \subset\left(E_{1}-E_{2}\right)^{\sim}=\tilde{E}_{1} \cup E_{2}$ and note that $Q=$ $\left(Q \cap E_{2}\right) \cup\left(Q-E_{2}\right)$. The $\varphi$-measurability of $E_{2}$ implies

$$
\begin{align*}
\varphi(P)+\varphi(Q) & =\varphi(P)+\varphi\left[\left(Q \cap E_{2}\right) \cup\left(Q-E_{2}\right)\right] \\
& =\varphi(P)+\varphi\left(Q \cap E_{2}\right)+\varphi\left(Q-E_{2}\right) \tag{4.3}
\end{align*}
$$

But $P \subset E_{1}, Q-E_{2} \subset \tilde{E}_{1}$, and the $\varphi$-measurability of $E_{1}$ imply

$$
\begin{align*}
\varphi(P)+\varphi\left(Q \cap E_{2}\right) & +\varphi\left(Q-E_{2}\right) \\
& =\varphi\left(Q \cap E_{2}\right)+\varphi\left[P \cup\left(Q-E_{2}\right)\right] \tag{4.4}
\end{align*}
$$

Also, $Q \cap E_{2} \subset E_{2}, P \cup\left(Q-E_{2}\right) \subset \tilde{E}_{2}$, and the $\varphi$-measurability of $E_{2}$ imply

$$
\begin{aligned}
\varphi\left(Q \cap E_{2}\right) & +\varphi\left[P \cup\left(Q-E_{2}\right)\right] \\
& =\varphi\left[\left(Q \cap E_{2}\right) \cup\left(P \cup\left(Q-E_{2}\right)\right)\right] \\
& =\varphi(Q \cup P)=\varphi(P \cup Q) .
\end{aligned}
$$

Hence, by (4.3), (4.4) and (4.5) we have

$$
\varphi(P)+\varphi(Q)=\varphi(P \cup Q)
$$

(iv) Let $S_{k}=\cup_{i=1}^{k} E_{i}$ and let $A$ be an arbitrary subset of $X$. We proceed by finite induction and first note that the result is obviously true for $k=1$. For $k>1$ assume $S_{k}$ is $\varphi$-measurable and that

$$
\begin{equation*}
\varphi(A) \geq \sum_{i=1}^{k} \varphi\left(A \cap E_{i}\right)+\varphi\left(A \cap \tilde{S}_{k}\right) \tag{4.6}
\end{equation*}
$$

for any set $A$. Then

$$
\begin{aligned}
& \varphi(A)=\varphi\left(A \cap E_{k+1}\right)+\varphi\left(A \cap \tilde{E}_{k+1}\right) \quad \text { because } E_{k+1} \text { is } \varphi \text {-measurable } \\
& =\varphi\left(A \cap E_{k+1}\right)+\varphi\left(A \cap \tilde{E}_{k+1} \cap S_{k}\right) \\
& +\varphi\left(A \cap \tilde{E}_{k+1} \cap \tilde{S}_{k}\right) \quad \text { because } S_{k} \text { is } \varphi \text {-measurable } \\
& =\varphi\left(A \cap E_{k+1}\right)+\varphi\left(A \cap S_{k}\right) \\
& +\varphi\left(A \cap \tilde{S}_{k+1}\right) \quad \text { because } S_{k} \subset \tilde{E}_{k+1} \\
& \geq \sum_{i=1}^{k+1} \varphi\left(A \cap E_{i}\right)+\varphi\left(A \cap \tilde{S}_{k+1}\right) \quad \text { use (4.6) with } A \text { replaced by } A \cap S_{k} \text {. }
\end{aligned}
$$

By the countable subadditivity of $\varphi$, this shows that

$$
\varphi(A) \geq \varphi\left(A \cap S_{k+1}\right)+\varphi\left(A \cap \tilde{S}_{k+1}\right)
$$

this, in turn, implies that $S_{k+1}$ is $\varphi$-measurable. Since we now know for any set $A \subset X$ and for all positive integers $k$ that

$$
\varphi(A) \geq \sum_{i=1}^{k} \varphi\left(A \cap E_{i}\right)+\varphi\left(A \cap \tilde{S}_{k}\right)
$$

and that $\tilde{S}_{k} \supset \tilde{S}$, we have

$$
\begin{align*}
\varphi(A) & \geq \sum_{i=1}^{\infty} \varphi\left(A \cap E_{i}\right)+\varphi(A \cap \tilde{S})  \tag{4.7}\\
& \geq \varphi(A \cap S)+\varphi(A \cap \tilde{S})
\end{align*}
$$

Again, the countable subadditivity of $\varphi$ was used to establish the last inequality. This implies that $S$ is $\varphi$-measurable, which establishes the first part of (iv). For the second part of (iv), note that the countable subadditivity of $\varphi$ yields

$$
\begin{aligned}
\varphi(A) & \leq \varphi(A \cap S)+\varphi(A \cap \tilde{S}) \\
& \leq \sum_{i=1}^{\infty} \varphi\left(A \cap E_{i}\right)+\varphi(A \cap \tilde{S})
\end{aligned}
$$

This, along with (4.7), establishes the last part of (iv).

The preceding result shows that $\varphi$-measurable sets are closed under the settheoretic operations of taking complements and countable disjoint unions. Of course, it would be preferable if they were closed under countable unions, and not merely countable disjoint unions. The proposition below addresses this issue. But first, we will prove a lemma that will be frequently used throughout. It states that the union of any countable family of sets can be written as the union of a countable family of disjoint sets.
4.7. Lemma. Let $\left\{E_{i}\right\}$ be a sequence of arbitrary sets. Then there exists a sequence of disjoint sets $\left\{A_{i}\right\}$ such that each $A_{i} \subset E_{i}$ and

$$
\bigcup_{i=1}^{\infty} E_{i}=\bigcup_{i=1}^{\infty} A_{i} .
$$

In case each $E_{i}$ is $\varphi$-measurable, so is $A_{i}$.
Proof. For each positive integer $j$, define $S_{j}=\cup_{i=1}^{j} E_{i}$. Note that

$$
\bigcup_{i=1}^{\infty} E_{i}=S_{1} \bigcup\left(\bigcup_{k=1}^{\infty}\left(S_{k+1} \backslash S_{k}\right)\right)
$$

Now take $A_{1}=S_{1}$ and $A_{i+1}=S_{i+1} \backslash S_{i}$ for all integers $i \geq 1$.
In case each $E_{i}$ is $\varphi$-measurable, the same is true for each $S_{j}$. Indeed, referring to Theorem 4.6 (iii), we see that $S_{2}$ is $\varphi$-measurable because $S_{2}=E_{2} \cup\left(E_{1} \backslash E_{2}\right)$ is the disjoint union of $\varphi$-measurable sets. Inductively, we see that $S_{j}=E_{j} \cup\left(S_{j-1} \backslash\right.$ $\left.E_{j}\right)$ is the disjoint union of $\varphi$-measurable sets and therefore the sets $A_{i}$ are also $\varphi$-measurable.
4.8. Theorem. If $\left\{E_{i}\right\}$ is a sequence of $\varphi$-measurable sets in $X$, then $\cup_{i=1}^{\infty} E_{i}$ and $\cap_{i=1}^{\infty} E_{i}$ are $\varphi$-measurable.

Proof. From the previous lemma, we have

$$
\bigcup_{i=1}^{\infty} E_{i}=\bigcup_{i=1}^{\infty} A_{i}
$$

where each $A_{i}$ is a $\varphi$-measurable subset of $E_{i}$ and where the sequence $\left\{A_{i}\right\}$ is disjoint. Thus, it follows immediately from Theorem 4.6 (iv), that $\cup_{i=1}^{\infty} E_{i}$ is $\varphi$ measurable.

To establish the second claim, note that

$$
X \backslash\left(\bigcap_{i=1}^{\infty} E_{i}\right)=\bigcup_{i=1}^{\infty} \tilde{E}_{i}
$$

The right side is $\varphi$-measurable in view of Theorem 4.6 (ii), (iii) and the first claim. By appealing again to Theorem 4.6 (iv), this concludes the proof.

Classes of sets that are closed under complementation and countable unions play an important role in measure theory and are therefore given a special name.
4.9. Definition. A nonempty collection $\Sigma$ of sets $E$ satisfying the following two conditions is called a $\sigma$-algebra:
(i) if $E \in \Sigma$, then $\tilde{E} \in \Sigma$,
(ii) $\cup_{i=1}^{\infty} E_{i} \in \Sigma$ provided each $E_{i} \in \Sigma$.

Note that it easily follows from the definition that a $\sigma$-algebra is closed under countable intersections and finite differences. Note also that the entire space and the empty set are elements of the $\sigma$-algebra since $\emptyset=E \cap \tilde{E} \in \Sigma$.
4.10. Definition. In a topological space, the elements of the smallest $\sigma$-algebra that contains all open sets are called Borel Sets. The term "smallest" is taken in the sense of inclusion and it is left as an exercise (Exercise 1, Section 4.2) to show that such a smallest $\sigma$-algebra, denoted as $\mathcal{B}$, exists.

The following is an immediate consequence of Theorem 4.6 and Theorem 4.8.
4.11. Corollary. If $\varphi$ is an outer measure on an arbitrary set $X$, then the class of $\varphi$-measurable sets forms a $\sigma$-algebra.

Next we state a result that exhibits the basic additivity and continuity properties of outer measure when restricted to its measurable sets. These properties follow almost immediately from Theorem 4.6.
4.12. Corollary. Suppose $\varphi$ is an outer measure on $X$ and $\left\{E_{i}\right\}$ a countable collection of $\varphi$-measurable sets.
(i) If $E_{1} \subset E_{2}$ with $\varphi\left(E_{1}\right)<\infty$, then

$$
\varphi\left(E_{2} \backslash E_{1}\right)=\varphi\left(E_{2}\right)-\varphi\left(E_{1}\right)
$$

(See Exercise 3, Section 4.1.)
(ii) (Countable additivity) If $\left\{E_{i}\right\}$ is a disjoint sequence of sets, then

$$
\varphi\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\sum_{i=1}^{\infty} \varphi\left(E_{i}\right)
$$

(iii) (Continuity from the left) If $\left\{E_{i}\right\}$ is an increasing sequence of sets, that is, if $E_{i} \subset E_{i+1}$ for each $i$, then

$$
\varphi\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\varphi\left(\lim _{i \rightarrow \infty} E_{i}\right)=\lim _{i \rightarrow \infty} \varphi\left(E_{i}\right) .
$$

(iv) (Continuity from the right) If $\left\{E_{i}\right\}$ is a decreasing sequence of sets, that is, if $E_{i} \supset E_{i+1}$ for each $i$, and if $\varphi\left(E_{i_{0}}\right)<\infty$ for some $i_{0}$, then

$$
\varphi\left(\bigcap_{i=1}^{\infty} E_{i}\right)=\varphi\left(\lim _{i \rightarrow \infty} E_{i}\right)=\lim _{i \rightarrow \infty} \varphi\left(E_{i}\right)
$$

(v) If $\left\{E_{i}\right\}$ is any sequence of $\varphi$-measurable sets, then

$$
\varphi\left(\liminf _{i \rightarrow \infty} E_{i}\right) \leq \liminf _{i \rightarrow \infty} \varphi\left(E_{i}\right)
$$

(vi) If

$$
\varphi\left(\bigcup_{i=i_{0}}^{\infty} E_{i}\right)<\infty
$$

for some positive integer $i_{0}$, then

$$
\varphi\left(\limsup _{i \rightarrow \infty} E_{i}\right) \geq \limsup _{i \rightarrow \infty} \varphi\left(E_{i}\right)
$$

Proof. We first observe that in view of Corollary 4.11, each of the sets that appears on the left side of (ii) through (vi) is $\varphi$-measurable. Consequently, all sets encountered in the proof will be $\varphi$-measurable.
(i): Observe that $\varphi\left(E_{2}\right)=\varphi\left(E_{2} \backslash E_{1}\right)+\varphi\left(E_{1}\right)$ since $E_{2} \backslash E_{1}$ is $\varphi$-measurable (Theorem 4.6 (iii)).
(ii): This is a restatement of Theorem 4.6 (iv).
(iii): We may assume that $\varphi\left(E_{i}\right)<\infty$ for each $i$, for otherwise the result follows from the monotonicity of $\varphi$. Since the sets $E_{1}, E_{2} \backslash E_{1}, \ldots, E_{i+1} \backslash E_{i}, \ldots$ are $\varphi$-measurable and disjoint, it follows that

$$
\lim _{i \rightarrow \infty} E_{i}=\bigcup_{i=1}^{\infty} E_{i}=E_{1} \cup\left[\bigcup_{i=1}^{\infty}\left(E_{i+1} \backslash E_{i}\right)\right]
$$

and therefore, from (iv) of Theorem 4.6, that

$$
\varphi\left(\lim _{i \rightarrow \infty} E_{i}\right)=\varphi\left(E_{1}\right)+\sum_{i=1}^{\infty} \varphi\left(E_{i+1} \backslash E_{i}\right)
$$

Since the sets $E_{i}$ and $E_{i+1} \backslash E_{i}$ are disjoint and $\varphi$-measurable, we have $\varphi\left(E_{i+1}\right)=$ $\varphi\left(E_{i+1} \backslash E_{i}\right)+\varphi\left(E_{i}\right)$. Therefore, because $\varphi\left(E_{i}\right)<\infty$ for each $i$, we have from (4.1)

$$
\begin{aligned}
\varphi\left(\lim _{i \rightarrow \infty} E_{i}\right) & =\varphi\left(E_{1}\right)+\sum_{i=1}^{\infty}\left[\varphi\left(E_{i+1}\right)-\varphi\left(E_{i}\right)\right] \\
& =\lim _{i \rightarrow \infty} \varphi\left(E_{i+1}\right)
\end{aligned}
$$

which proves (iii).
(iv): By replacing $E_{i}$ with $E_{i} \cap E_{i_{0}}$ if necessary, we may assume that $\varphi\left(E_{1}\right)<\infty$. Since $\left\{E_{i}\right\}$ is decreasing, the sequence $\left\{E_{1} \backslash E_{i}\right\}$ is increasing and therefore (iii) implies

$$
\begin{align*}
\varphi\left(\bigcup_{i=1}^{\infty}\left(E_{1} \backslash E_{i}\right)\right) & =\lim _{i \rightarrow \infty} \varphi\left(E_{1} \backslash E_{i}\right)  \tag{4.8}\\
& =\varphi\left(E_{1}\right)-\lim _{i \rightarrow \infty} \varphi\left(E_{i}\right)
\end{align*}
$$

It is easy to verify that

$$
\bigcup_{i=1}^{\infty}\left(E_{1} \backslash E_{i}\right)=E_{1} \backslash \bigcap_{i=1}^{\infty} E_{i}
$$

and therefore, from (i) of Corollary 4.12, we have

$$
\varphi\left(\bigcup_{i=1}^{\infty}\left(E_{1} \backslash E_{i}\right)\right)=\varphi\left(E_{1}\right)-\varphi\left(\bigcap_{i=1}^{\infty} E_{i}\right)
$$

which, along with (4.8), yields

$$
\varphi\left(E_{1}\right)-\varphi\left(\bigcap_{i=1}^{\infty} E_{i}\right)=\varphi\left(E_{1}\right)-\lim _{i \rightarrow \infty} \varphi\left(E_{i}\right)
$$

The fact that $\varphi\left(E_{1}\right)<\infty$ allows us to conclude (iv).
(v): Let $A_{j}=\cap_{i=j}^{\infty} E_{i}$ for $j=1,2, \ldots$. Then, $A_{j}$ is an increasing sequence of $\varphi$ measurable sets with the property that $\lim _{j \rightarrow \infty} A_{j}=\liminf _{i \rightarrow \infty} E_{i}$ and therefore, by (iii),

$$
\varphi\left(\liminf _{i \rightarrow \infty} E_{i}\right)=\lim _{j \rightarrow \infty} \varphi\left(A_{j}\right)
$$

But since $A_{j} \subset E_{j}$, it follows that

$$
\lim _{j \rightarrow \infty} \varphi\left(A_{j}\right) \leq \liminf _{j \rightarrow \infty} \varphi\left(E_{j}\right)
$$

thus establishing (v).
The proof of (vi) is similar to that of (v) and is left as Exercise 1, Section 4.1.
4.13. Remark. We mentioned earlier that one of our major concerns is to determine whether there is a rich supply of measurable sets for a given outer measure $\varphi$. Although we have learned that the class of measurable sets constitutes a $\sigma$ algebra, this is not sufficient to guarantee that the measurable sets exist in great numbers. For example, suppose that $X$ is an arbitrary set and $\varphi$ is defined on $X$ as $\varphi(E)=1$ whenever $E \subset X$ is nonempty while $\varphi(\emptyset)=0$. Then it is easy to verify that $X$ and $\emptyset$ are the only $\varphi$-measurable sets. In order to overcome this difficulty, it is necessary to impose an additivity condition on $\varphi$. This will be developed in the following section.

We will need the following definitions:
4.14. Definitions. An outer measure $\varphi$ on a topological space $X$ is called a Borel outer measure if all Borel sets are $\varphi$-measurable. A Borel outer measure is finite if $\varphi(X)$ is finite.

An outer measure $\varphi$ on a set $X$ is called regular if for each $A \subset X$ there exists a $\varphi$-measurable set $B \supset A$ such that $\varphi(B)=\varphi(A)$. A Borel regular outer measure is a Borel outer measure such that for each $A \subset X$, there exists a Borel set $B$ such that $\varphi(B)=\varphi(A)$ (see Theorem 4.52 and Corollary 4.56 for regularity
properties of Borel outer measures). A Radon outer measure is a Borel regular outer measure that is finite on compact sets.

We began this section with the concept of an outer measure on an arbitrary set $X$ and proved that the family of $\varphi$-measurable sets forms a $\sigma$-algebra. We will see in Section 4.9 that the restriction of $\varphi$ to this $\sigma$-algebra generates a measure space (see Definition 4.47). In the next few sections we will introduce important examples of outer measures, which in turn will provide measure spaces that appear in many areas of mathematics.

## Exercises for Section 4.1

1. Prove (vi) of Corollary 4.12.
2. In example (iv), let $\varphi_{\varepsilon}(A):=\varphi(A)$ to denote the dependence on $\varepsilon$, and define

$$
\psi(A):=\lim _{\varepsilon \rightarrow 0} \varphi_{\varepsilon}(A)
$$

What is $\psi(A)$ and what are the corresponding $\psi$-measurable sets?
3. In (i) of Corollary 4.12, it was shown that $\varphi\left(E_{2} \backslash E_{1}\right)=\varphi\left(E_{2}\right)-\varphi\left(E_{1}\right)$ provided $E_{1} \subset E_{2}$ are $\varphi$-measurable with $\varphi\left(E_{1}\right)<\infty$. Prove this result still remains true if $E_{2}$ is not assumed to be $\varphi$-measurable.

### 4.2. Carathéodory Outer Measure

In the previous section, we considered an outer measure $\varphi$ on an arbitrary set $X$. We now restrict our attention to a metric space $X$ and impose a further condition (an additivity condition) on the outer measure. This will allow us to conclude that all closed sets are measurable.
4.15. Definition. An outer measure $\varphi$ defined on a metric space $(X, \rho)$ is called a Carathéodory outer measure if

$$
\begin{equation*}
\varphi(A \cup B)=\varphi(A)+\varphi(B) \tag{4.9}
\end{equation*}
$$

whenever $A, B$ are arbitrary subsets of $X$ with $\boldsymbol{d}(A, B)>0$. The notation $\boldsymbol{d}(A, B)$ denotes the distance between the sets $A$ and $B$ and is defined by

$$
\boldsymbol{d}(A, B):=\inf \{\rho(a, b): a \in A, b \in B\}
$$

4.16. Theorem. If $\varphi$ is a Carathéodory outer measure on a metric space $X$, then all closed sets are $\varphi$-measurable.

Proof. We will verify the condition in Definition 4.3 whenever $C$ is a closed set. Because $\varphi$ is subadditive, it suffices to show

$$
\begin{equation*}
\varphi(A) \geq \varphi(A \cap C)+\varphi(A \backslash C) \tag{4.10}
\end{equation*}
$$

whenever $A \subset X$. In order to prove (4.10), consider $A \subset X$ with $\varphi(A)<\infty$ and for each positive integer $i$, let $C_{i}=\{x: \boldsymbol{d}(x, C) \leq 1 / i\}$. Note that

$$
\boldsymbol{d}\left(A \backslash C_{i}, A \cap C\right) \geq \frac{1}{i}>0
$$

Since $A \supset\left(A \backslash C_{i}\right) \cup(A \backslash C)$, (4.15) implies

$$
\begin{equation*}
\varphi(A) \geq \varphi\left(\left(A \backslash C_{i}\right) \cup(A \cap C)\right)=\varphi\left(A \backslash C_{i}\right)+\varphi(A \cap C) \tag{4.11}
\end{equation*}
$$

Because of this inequality, the proof of (4.10) will be concluded if we can show that

$$
\begin{equation*}
\lim _{i \rightarrow \infty} \varphi\left(A \backslash C_{i}\right)=\varphi(A \backslash C) \tag{4.12}
\end{equation*}
$$

For each positive integer $i$, let

$$
T_{i}=A \cap\left\{x: \frac{1}{i+1}<\boldsymbol{d}(x, C) \leq \frac{1}{i}\right\}
$$

and note that since $C$ is closed, $x \notin C$ if and only if $\boldsymbol{d}(x, C)>0$ and therefore that

$$
\begin{equation*}
A \backslash C=\left(A \backslash C_{j}\right) \cup\left(\bigcup_{i=j}^{\infty} T_{i}\right) \tag{4.13}
\end{equation*}
$$

for each positive integer $j$. This, in turn, implies

$$
\begin{equation*}
\varphi(A \backslash C) \leq \varphi\left(A \backslash C_{j}\right)+\sum_{i=j}^{\infty} \varphi\left(T_{i}\right) \tag{4.14}
\end{equation*}
$$

We now note that

$$
\begin{equation*}
\sum_{i=1}^{\infty} \varphi\left(T_{i}\right)<\infty \tag{4.15}
\end{equation*}
$$

To establish (4.15), first observe that $\boldsymbol{d}\left(T_{i}, T_{j}\right)>0$ if $|i-j| \geq 2$. Thus, we obtain from (4.9) that for each positive integer $m$,

$$
\begin{aligned}
\sum_{i=1}^{m} \varphi\left(T_{2 i}\right) & =\varphi\left(\bigcup_{i=1}^{m} T_{2 i}\right) \leq \varphi(A)<\infty \\
\sum_{i=1}^{m} \varphi\left(T_{2 i-1}\right) & =\varphi\left(\bigcup_{i=1}^{m} T_{2 i-1}\right) \leq \varphi(A)<\infty
\end{aligned}
$$

From (4.14) and since $A \backslash C_{j} \subset A \backslash C$ and $\sum_{i=1}^{\infty} \varphi\left(T_{i}\right)<\infty$ we have

$$
\varphi(A \backslash C)-\sum_{i=j}^{\infty} \varphi\left(T_{i}\right) \leq \varphi\left(A \backslash C_{j}\right) \leq \varphi(A \backslash C)
$$

Hence, by letting $j \rightarrow \infty$ an using $\lim _{j \rightarrow \infty} \sum_{i=j}^{\infty} \varphi\left(T_{i}\right)=0$ we obtain the desired conclusion.

The following proposition provides a useful description of the Borel sets.
4.17. Theorem. Suppose $\mathcal{F}$ is a family of subsets of a topological space $X$ that contains all open and all closed subsets of $X$. Suppose also that $\mathcal{F}$ is closed under countable unions and countable intersections. Then $\mathcal{F}$ contains all Borel sets; that is, $\mathcal{B} \subset \mathcal{F}$.

Proof. Let

$$
\mathcal{H}=\mathcal{F} \cap\{A: \tilde{A} \in \mathcal{F}\}
$$

Observe that $\mathcal{H}$ contains all closed sets. Moreover, it is easily seen that $\mathcal{H}$ is closed under complementation and countable unions. Thus, $\mathcal{H}$ is a $\sigma$-algebra that contains the open sets and therefore contains all Borel sets.

As a direct result of Corollary 4.11 and Theorem 4.16 we have the main result of this section.
4.18. Theorem. If $\varphi$ is a Carathéodory outer measure on a metric space $X$, then the Borel sets of $X$ are $\varphi$-measurable.

In case $X=\mathbb{R}^{n}$, it follows that the cardinality of the Borel sets is at least as great as that of the closed sets. Since the Borel sets contain all singletons of $\mathbb{R}^{n}$, their cardinality is at least $c$. We thus have shown that not only do the $\varphi$ measurable sets have nice additivity properties (they form a $\sigma$-algebra), but in addition, there is a plentiful supply of them in case $\varphi$ is an Carathéodory outer measure on $\mathbb{R}^{n}$. Thus, the difficulty that arises from the example in Remark 4.13 is avoided. In the next section we discuss a concrete illustration of such a measure.

## Exercises for Section 4.2

1. Prove that, in any topological space $X$, there exists a smallest $\sigma$-algebra that contains all open sets in $X$. That is, prove there is a $\sigma$-algebra $\Sigma$ that contains all open sets and has the property that if $\Sigma_{1}$ is another $\sigma$-algebra containing all open sets, then $\Sigma \subset \Sigma_{1}$. In particular, for $X=\mathbb{R}^{n}$, note that there is a smallest $\sigma$-algebra that contains all the closed sets in $\mathbb{R}^{n}$.
2. In a topological space $X$ the family of Borel sets, $\mathcal{B}$, is by definition, the $\sigma$-algebra generated by the closed sets. The method below is another way of describing the Borel sets using transfinite induction. You are to fill in the necessary steps: (a) For an arbitrary family $\mathcal{F}$ of sets, let

$$
\mathcal{F}^{*}=\left\{\bigcup_{k=1}^{\infty} E_{k}: \text { where either } E_{i} \in \mathcal{F} \text { or } \widetilde{E}_{i} \in \mathcal{F} \text { for all } i \in \mathbb{N}\right\}
$$

Let $\Omega$ denote the smallest uncountable ordinal. We will use transfinite induction to define a family $\mathcal{E}_{\alpha}$ for each $\alpha<\Omega$.
(b) Let $\mathcal{E}_{0}:=$ all closed sets $:=\mathcal{K}$. Now choose $\alpha<\Omega$ and assume $\mathcal{E}_{\beta}$ has been defined for each $\beta$ such that $0 \leq \beta<\alpha$. Define

$$
\mathcal{E}_{\alpha}:=\left(\bigcup_{0 \leq \beta<\alpha} \mathcal{E}_{\beta}\right)^{*}
$$

and define

$$
\mathcal{A}:=\bigcup_{0 \leq \alpha<\Omega} \mathcal{E}_{\alpha}
$$

(c) Show that each $\mathcal{E}_{\alpha} \in \mathcal{B}$.
(d) Show that $\mathcal{A} \subset \mathcal{B}$
(e) Now show that $\mathcal{A}$ is a $\sigma$-algebra to conclude that $\mathcal{A}=\mathcal{B}$.
(i) Show that $\emptyset, X \in \mathcal{A}$.
(ii) Let $A \in \mathcal{A} \Longrightarrow A \in \mathcal{E}_{\alpha}$ for some $\alpha<\Omega$. Show that this $\Longrightarrow$ $\widetilde{A} \in \mathcal{E}_{\alpha}^{*} \subset \mathcal{E}_{\beta}$ for every $\beta>\alpha$.
(iii) Conclude that $\widetilde{A} \in \mathcal{A}$ and thus conclude that $\mathcal{A}$ is closed under complementation.
(iv) Now let $\left\{A_{k}\right\}$ be a sequence in $\mathcal{A}$. Show that $\bigcup_{k=1}^{\infty} A_{k} \in \mathcal{A}$. (Hint: Each $A_{k} \in \mathcal{A}_{\alpha_{k}}$ for some $\alpha_{k}<\Omega$. We know there is $\beta<\Omega$ such that $\beta>\alpha_{k}$ for each $k \in \mathbb{N}$ ).
3. Prove that the set function $\mu$ defined in (4.40) is an outer measure whose measurable sets include all open sets.
4. Prove that the set function $\psi$ defined in (4.45) is an outer measure on $X$.
5. An outer measure $\varphi$ on a space $X$ is called $\sigma$-finite if there exists a countable number of sets $A_{i}$ with $\varphi\left(A_{i}\right)<\infty$ such that $X \subset \cup_{i=1}^{\infty} A_{i}$. Assuming $\varphi$ is a $\sigma$-finite Borel regular outer measure on a metric space $X$, prove that $E \subset X$ is $\varphi$-measurable if and only if there exists an $F_{\sigma}$ set $F \subset E$ such that $\varphi(E \backslash F)=0$.
6. Let $\varphi$ be an outer measure on a space $X$. Suppose $A \subset X$ is an arbitrary set with $\varphi(A)<\infty$ and such that there exists a $\varphi$-measurable set $E \supset A$ with $\varphi(E)=\varphi(A)$. Prove that $\varphi(A \cap B)=\varphi(E \cap B)$ for every $\varphi$-measurable set $B$.
7. In $\mathbb{R}^{2}$, find two disjoint closed sets $A$ and $B$ such that $\boldsymbol{d}(A, B)=0$. Show this is not possible if one of the sets is compact.
8. Let $\varphi$ be an outer Carathéodory measure on $\mathbb{R}$ and let $f(x):=\varphi\left(I_{x}\right)$ where $I_{x}$ is an open interval of fixed length centered at $x$. Prove that $f$ is lower semicontinuous. What can you say about $f$ if $I_{x}$ is taken as a closed interval? Prove the analogous result in $\mathbb{R}^{n}$; that is, let $f(x):=\varphi(B(x, a))$ where $B(x, a)$ is the open ball with fixed radius $a$ and centered at $x$.
9. In a metric space $X$, prove that dist $(A, B)=\boldsymbol{d}(\bar{A}, \bar{B})$ for arbitrary sets $A, B \in$ $X$.
10. Let $A$ be a non-Borel subset of $\mathbb{R}^{n}$ and define for each subset $E$,

$$
\varphi(E)= \begin{cases}0 & \text { if } E \subset A \\ \infty & \text { if } E \backslash A \neq \emptyset\end{cases}
$$

Prove that $\varphi$ is an outer measure that is not Borel regular.
11. Let $\mathcal{M}$ denote the class of $\varphi$-measurable sets of an outer measure $\varphi$ defined on a set $X$. If $\varphi(X)<\infty$, prove that the family

$$
\mathcal{F}:=\{A \in \mathcal{M}: \varphi(A)>0\}
$$

is at most countable.

### 4.3. Lebesgue Measure

Lebesgue measure on $\mathbb{R}^{n}$ is perhaps the most important example of a Carathéodory outer measure. We will investigate the properties of this measure and show, among other things, that it agrees with the primitive notion of volume on sets such as $n$-dimensional "intervals."

For the purpose of defining Lebesgue outer measure on $\mathbb{R}^{n}$, we consider closed $n$-dimensional intervals

$$
\begin{equation*}
I=\left\{x: a_{i} \leq x_{i} \leq b_{i}, i=1,2, \ldots, n\right\} \tag{4.16}
\end{equation*}
$$

and their volumes

$$
\begin{equation*}
v(I)=\prod_{i=1}^{n}\left(b_{i}-a_{i}\right) \tag{4.17}
\end{equation*}
$$

With $I_{1}=\left[a_{1}, b_{1}\right], I_{2}=\left[a_{2}, b_{2}\right], \ldots, I_{n}=\left[a_{n}, b_{n}\right]$, we have

$$
I=I_{1} \times I_{2} \times \cdots \times I_{n}
$$

Notice that $n$-dimensional intervals have their edges parallel to the coordinate axes of $\mathbb{R}^{n}$. When no confusion arises, we shall simply say "interval" rather than " $n$ dimensional interval."

In preparation for the development of Lebesgue measure, we state two elementary propositions concerning intervals whose proofs will be omitted.
4.19. Theorem. Suppose each edge $I_{k}=\left[a_{k}, b_{k}\right]$ of an $n$-dimensional interval $I$ is partitioned into $\alpha_{k}$ subintervals. The products of these intervals produce a partition of $I$ into $\beta:=\alpha_{1} \cdot \alpha_{2} \cdots \alpha_{n}$ subintervals $I_{i}$ and

$$
v(I)=\sum_{i=1}^{\beta} v\left(I_{i}\right)
$$

4.20. Theorem. For each interval $I$ and each $\varepsilon>0$, there exists an interval $J$ whose interior contains I and

$$
v(J)<v(I)+\varepsilon
$$

4.21. Definition. The Lebesgue outer measure of an arbitrary set $E \subset \mathbb{R}^{n}$, denoted by $\lambda^{*}(E)$, is defined by

$$
\lambda^{*}(E)=\inf \left\{\sum_{k=1}^{\infty} v\left(I_{k}\right)\right\}
$$

where the infimum is taken over all countable collections of closed intervals $I_{k}$ such that

$$
E \subset \bigcup_{k=1}^{\infty} I_{k} .
$$

It may be necessary at times to emphasize the dimension of the Euclidean space in which Lebesgue outer measure is defined. When clarification is needed, we will write $\lambda_{n}^{*}(E)$ in place of $\lambda^{*}(E)$.

Our next result shows that Lebesgue outer measure is an extension of volume.
4.22. Theorem. For a closed interval $I \subset \mathbb{R}^{n}, \lambda^{*}(I)=v(I)$.

Proof. The inequality $\lambda^{*}(I) \leq v(I)$ holds since $\mathcal{S}$ consisting of $I$ alone can be taken as one of the admissible competitors in Definition 4.21.

To prove the opposite inequality, choose $\varepsilon>0$ and let $\left\{I_{k}\right\}_{k=1}^{\infty}$ be a sequence of closed intervals such that

$$
\begin{equation*}
I \subset \bigcup_{k=1}^{\infty} I_{k} \quad \text { and } \quad \sum_{k=1}^{\infty} v\left(I_{k}\right)<\lambda^{*}(I)+\varepsilon \tag{4.18}
\end{equation*}
$$

For each $k$, refer to Proposition 4.20 to obtain an interval $J_{k}$ whose interior contains $I_{k}$ and

$$
v\left(J_{k}\right) \leq v\left(I_{k}\right)+\frac{\varepsilon}{2^{k}}
$$

We therefore have

$$
\sum_{k=1}^{\infty} v\left(J_{k}\right) \leq \sum_{k=1}^{\infty} v\left(I_{k}\right)+\varepsilon
$$

Let $\mathcal{F}=\left\{\right.$ interior $\left.\left(J_{k}\right): k \in \mathbb{N}\right\}$ and observe that $\mathcal{F}$ is an open cover of the compact set $I$. Let $\eta$ be the Lebesgue number for $\mathcal{F}$ (see Exercise 4, Section 3.5). By Proposition 4.19, there is a partition of $I$ into finitely many subintervals, $K_{1}, K_{2}, \ldots, K_{m}$, each with diameter less than $\eta$ and having the property

$$
I=\bigcup_{i=1}^{m} K_{i} \quad \text { and } \quad v(I)=\sum_{i=1}^{m} v\left(K_{i}\right)
$$

Each $K_{i}$ is contained in the interior of some $J_{k}$, say $J_{k_{i}}$, although more than one $K_{i}$ may belong to the same $J_{k_{i}}$. Thus, if $N_{m}$ denotes the smallest number of the $J_{k_{i}}$ 's that contain the $K_{i}$ 's, we have $N_{m} \leq m$ and

$$
v(I)=\sum_{i=1}^{m} v\left(K_{i}\right) \leq \sum_{i=1}^{N_{m}} v\left(J_{k_{i}}\right) \leq \sum_{k=1}^{\infty} v\left(J_{k}\right) \leq \sum_{k=1}^{\infty} v\left(I_{k}\right)+\varepsilon
$$

From this and (4.18) it follows that

$$
v(I) \leq \lambda^{*}(I)+2 \varepsilon
$$

which yields the desired result since $\varepsilon$ is arbitrary.
We will now show that Lebesgue outer measure is a Carathéodory outer measure as defined in Definition 4.15. Once we have established this result, we then will be able to apply the important results established in Section 4.2, such as Theorem 4.18, to Lebesgue outer measure.
4.23. ThEOREM. Lebesgue outer measure, $\lambda^{*}$, defined on $\mathbb{R}^{n}$ is a Carathéodory outer measure.

Proof. We first verify that $\lambda^{*}$ is an outer measure. The first three conditions of Definition 4.1 are immediate, so we proceed with the proof of condition (iv). Let $\left\{A_{i}\right\}$ be a countable collection of arbitrary sets in $\mathbb{R}^{n}$ and let $A=\cup_{i=1}^{\infty} A_{i}$. We may as well assume that $\lambda^{*}\left(A_{i}\right)<\infty$ for $i=1,2, \ldots$, for otherwise the conclusion is obvious. Choose $\varepsilon>0$. For each $i$, the definition of Lebesgue outer measure implies that there exists a countable family of closed intervals, $\left\{I_{j}^{(i)}\right\}_{j=1}^{\infty}$, such that

$$
\begin{equation*}
A_{i} \subset \bigcup_{j=1}^{\infty} I_{j}^{(i)} \tag{4.19}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{j=1}^{\infty} v\left(I_{j}^{(i)}\right)<\lambda^{*}\left(A_{i}\right)+\frac{\varepsilon}{2^{i}} \tag{4.20}
\end{equation*}
$$

Now $A \subset \bigcup_{i, j=1}^{\infty} I_{j}^{(i)}$ and therefore

$$
\begin{aligned}
\lambda^{*}(A) & \leq \sum_{i, j=1}^{\infty} v\left(I_{j}^{(i)}\right)=\sum_{i=1}^{\infty} \sum_{j=1}^{\infty} v\left(I_{j}^{(i)}\right) \\
& \leq \sum_{i=1}^{\infty}\left(\lambda^{*}\left(A_{i}\right)+\frac{\varepsilon}{2^{i}}\right)=\sum_{i=1}^{\infty} \lambda^{*}\left(A_{i}\right)+\varepsilon
\end{aligned}
$$

Since $\varepsilon>0$ is arbitrary, the countable subadditivity of $\lambda^{*}$ is established.

Finally, we verify (4.15) of Definition 4.15 . Let $A$ and $B$ be arbitrary sets with $\boldsymbol{d}(A, B)>0$. From what has just been proved, we know that $\lambda^{*}(A \cup B) \leq \lambda^{*}(A)+$ $\lambda^{*}(B)$. To prove the opposite inequality, choose $\varepsilon>0$ and, from the definition of Lebesgue outer measure, select closed intervals $\left\{I_{k}\right\}$ whose union contains $A \cup B$ such that

$$
\sum_{k=1}^{\infty} v\left(I_{k}\right) \leq \lambda^{*}(A \cup B)+\varepsilon
$$

By subdividing each interval $I_{k}$ into smaller intervals if necessary, we may assume that the diameter of each $I_{k}$ is less than $\boldsymbol{d}(A, B)$. Thus, the family $\left\{I_{k}\right\}$ consists of two subfamilies, $\left\{I_{k}^{\prime}\right\}$ and $\left\{I_{k}^{\prime \prime}\right\}$, where the elements of the first have nonempty intersections with $A$ while the elements of the second have nonempty intersections with $B$. Consequently,

$$
\lambda^{*}(A)+\lambda^{*}(B) \leq \sum_{k=1}^{\infty} v\left(I_{k}^{\prime}\right)+\sum_{k=1}^{\infty} v\left(I_{k}^{\prime \prime}\right)=\sum_{k=1}^{\infty} v\left(I_{k}\right) \leq \lambda^{*}(A \cup B)+\varepsilon
$$

Since $\varepsilon>0$ is arbitrary, this shows that

$$
\lambda^{*}(A)+\lambda^{*}(B) \leq \lambda^{*}(A \cup B)
$$

which completes the proof.
4.24. Remark. We will henceforth refer to $\lambda^{*}$-measurable sets as Lebesgue measurable sets. Now that we know that Lebesgue outer measure is a Carathéodory outer measure, it follows from Theorem 4.18 that all Borel sets in $\mathbb{R}^{n}$ are Lebesgue measurable. In particular, each open set and each closed set is Lebesgue measurable. We will denote by $\lambda$ the set function obtained by restricting $\lambda^{*}$ to the family of Lebesgue measurable sets. Thus, whenever $E$ is a Lebesgue measurable set, we have by definition $\lambda(E)=\lambda^{*}(E)$. $\lambda$ is called Lebesgue measure. Note that the additivity and continuity properties established in Corollary 4.12 apply to Lebesgue measure.

In view of Theorem 4.22 and the continuity properties of Lebesgue measure, it is possible to show that the Lebesgue measure of elementary geometric figures in $\mathbb{R}^{n}$ agrees with the notion of volume. For example, suppose that $J$ is an open interval in $\mathbb{R}^{n}$, that is, suppose $J$ is the product of open 1-dimensional intervals. It is easily seen that $\lambda(J)$ equals the product of lengths of these intervals because $J$ can be written as the union of an increasing sequence $\left\{I_{k}\right\}$ of closed intervals. Then

$$
\lambda(J)=\lim _{k \rightarrow \infty} \lambda\left(I_{k}\right)=\lim _{k \rightarrow \infty} \operatorname{vol}\left(I_{k}\right)=\operatorname{vol}(J)
$$

Next, we give several characterizations of Lebesgue measurable sets. We recall Definition 3.46 , in which the concepts of $G_{\delta}$ and $F_{\sigma}$ sets are introduced.
4.25. Theorem. The following five conditions are equivalent for Lebesgue outer measure, $\lambda^{*}$, on $\mathbb{R}^{n}$.
(i) $E \subset \mathbb{R}^{n}$ is $\lambda^{*}$-measurable.
(ii) For each $\varepsilon>0$, there is an open set $U \supset E$ such that $\lambda^{*}(U \backslash E)<\varepsilon$.
(iii) There is a $G_{\delta}$ set $U \supset E$ such that $\lambda^{*}(U \backslash E)=0$.
(iv) For each $\varepsilon>0$, there is a closed set $F \subset E$ such that $\lambda^{*}(E \backslash F)<\varepsilon$.
(v) There is a $F_{\sigma}$ set $F \subset E$ such that $\lambda^{*}(E \backslash F)=0$.

Proof. (i) $\Rightarrow$ (ii). We first assume that $\lambda(E)<\infty$. For arbitrary $\varepsilon>0$, the definition of Lebesgue outer measure implies the existence of closed $n$-dimensional intervals $I_{k}$ whose union contains $E$ and such that

$$
\sum_{k=1}^{\infty} v\left(I_{k}\right)<\lambda^{*}(E)+\frac{\varepsilon}{2}
$$

Now, for each $k$, let $I_{k}^{\prime}$ be an open interval containing $I_{k}$ such that $v\left(I_{k}^{\prime}\right)<v\left(I_{k}\right)+$ $\varepsilon / 2^{k+1}$. Then, defining $U=\cup_{k=1}^{\infty} I_{k}^{\prime}$, we have that $U$ is open and from (4.3), that

$$
\lambda(U) \leq \sum_{k=1}^{\infty} v\left(I_{k}^{\prime}\right)<\sum_{k=1}^{\infty} v\left(I_{k}\right)+\varepsilon / 2<\lambda^{*}(E)+\varepsilon .
$$

Thus, $\lambda(U)<\lambda^{*}(E)+\varepsilon$, and since $E$ is a Lebesgue measurable set of finite measure, we may appeal to Corollary 4.12 (i) to conclude that

$$
\lambda(U \backslash E)=\lambda^{*}(U \backslash E)=\lambda^{*}(U)-\lambda^{*}(E)<\varepsilon
$$

In case $\lambda(E)=\infty$, for each positive integer $i$, let $E_{i}$ denote $E \cap B(i)$ where $B(i)$ is the open ball of radius $i$ centered at the origin. Then $E_{i}$ is a Lebesgue measurable set of finite measure and thus we may apply the previous step to find an open set $U_{i} \supset E_{i}$ such that $\lambda\left(U_{i} \backslash E_{i}\right)<\varepsilon / 2^{i}$. Let $U=\cup_{i=1}^{\infty} U_{i}$ and observe that $U \backslash E \subset \cup_{i=1}^{\infty}\left(U_{i} \backslash E_{i}\right)$. Now use the subadditivity of $\lambda$ to conclude that

$$
\lambda(U \backslash E) \leq \sum_{i=1}^{\infty} \lambda\left(U_{i} \backslash E_{i}\right)<\sum_{i=1}^{\infty} \frac{\varepsilon}{2^{i}}=\varepsilon
$$

which establishes the implication (i) $\Rightarrow$ (ii).
(ii) $\Rightarrow$ (iii). For each positive integer $i$, let $U_{i}$ denote an open set with the property that $U_{i} \supset E$ and $\lambda^{*}\left(U_{i} \backslash E\right)<1 / i$. If we define $U=\cap_{i=1}^{\infty} U_{i}$, then

$$
\lambda^{*}(U \backslash E)=\lambda^{*}\left[\bigcap_{i=1}^{\infty}\left(U_{i} \backslash E\right)\right] \leq \lim _{i \rightarrow \infty} \frac{1}{i}=0
$$

(iii) $\Rightarrow$ (i). This is obvious since both $U$ and $(U \backslash E)$ are Lebesgue measurable sets with $E=U \backslash(U \backslash E)$.
(i) $\Rightarrow$ (iv). Assume that $E$ is a measurable set and thus that $\tilde{E}$ is measurable. We know that (ii) is equivalent to (i) and thus, given $\varepsilon>0$, there is an open set $U \supset \tilde{E}$ such that $\lambda(U \backslash \tilde{E})<\varepsilon$. Note that

$$
E \backslash \tilde{U}=E \cap U=U \backslash \tilde{E}
$$

Since $\tilde{U}$ is closed, $\tilde{U} \subset E$ and $\lambda(E \backslash \tilde{U})<\varepsilon$, we see that (iv) holds with $F=\tilde{U}$.
The proofs of (iv) $\Rightarrow$ (v) and (v) $\Rightarrow$ (i) are analogous to those of (ii) $\Rightarrow$ (iii) and (iii) $\Rightarrow$ (i), respectively.
4.26. REMARK. The above proof is direct and uses only the definition of Lebesgue measure to establish the various regularity properties. However, another proof, which is not so long but is perhaps less transparent, proceeds as follows. Using only the definition of Lebesgue measure, it can be shown that for any set $A \subset \mathbb{R}^{n}$, there is a $G_{\delta}$ set $G \supset A$ such that $\lambda(G)=\lambda^{*}(A)$ (see Exercise 9 , Section 4.3). Since $\lambda^{*}$ is a Carathéodory outer measure (Theorem 4.23), its measurable sets contain the Borel sets (Theorem 4.18). Consequently, $\lambda^{*}$ is a Borel regular outer measure and thus, we may appeal to Corollary 4.56 below to conclude that assertions (ii) and (iv) of Theorem 4.25 hold for any Lebesgue measurable set. The remaining properties follow easily from these two.

## Exercises for Section 4.3

1. With $\lambda_{t}$ defined by

$$
\lambda_{t}(E)=\lambda(|t| E)
$$

prove that $\lambda_{t}(N)=0$ whenever $\lambda(N)=0$.
2. Let $I, I_{1}, I_{2}, \ldots, I_{k}$ be intervals in $\mathbb{R}$ such that $I \subset \cup_{i=1}^{k} I_{i}$. Prove that

$$
v(I) \leq \sum_{i=1}^{k} v\left(I_{i}\right)
$$

where $v(I)$ denotes the length of the interval $I$.
3. Complete the proofs of (iv) $\Rightarrow$ (v) and (v) $\Rightarrow$ (i) in Theorem 4.25.
4. Let $E \subset \mathbb{R}$ and for each real number $t$, let $E+t=\{x+t: x \in E\}$. Prove that $\lambda^{*}(E)=\lambda^{*}(E+t)$. From this show that if $E$ is Lebesgue measurable, then so is $E+t$.
5. Prove that Lebesgue measure on $\mathbb{R}^{n}$ is independent of the choice of coordinate system. That is, prove that Lebesgue outer measure is invariant under rigid motions in $\mathbb{R}^{n}$.
6. Let $P$ denote an arbitrary $n-1$-dimensional hyperplane in $\mathbb{R}^{n}$. Prove that $\lambda(P)=0$.
7. In this problem, we want to show that any Lebesgue measurable subset of $\mathbb{R}$ must be "densely populated" in some interval. Thus, let $E \subset \mathbb{R}$ be a Lebesgue measurable set, $\lambda(E)>0$. For each $\varepsilon>0$, show that there exists an interval $I$ such that $\frac{\lambda(E \cap I)}{\lambda(I)}>1-\varepsilon$.
8. Suppose $E \subset \mathbb{R}^{n}, \lambda^{*}(E)<\infty$, is an arbitrary set with the property that there exists an $F_{\sigma}$-set $F \subset E$ with $\lambda(F)=\lambda^{*}(E)$. Prove that $E$ is a Lebesgue measurable set.
9. Prove that any arbitrary set $A \subset \mathbb{R}^{n}$ is contained within a $G_{\delta}$-set $G$ with the property $\lambda(G)=\lambda^{*}(A)$.
10. Let $\left\{E_{k}\right\}$ be a sequence of Lebesgue measurable sets contained in a compact set $K \subset \mathbb{R}^{n}$. Assume for some $\varepsilon>0$, that $\lambda\left(E_{k}\right)>\varepsilon$ for all $k$. Prove that there is some point that belongs to infinitely many $E_{k}$ 's.
11. Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a Lipschitz map. Prove that if $\lambda(E)=0$, then $\lambda(T(E))=0$.

### 4.4. The Cantor Set


#### Abstract

The Cantor set construction discussed in this section provides a method of generating a wide variety of important, and often unexpected, examples in real analysis. One of our main interests here is to show how the Cantor set exhibits the disparities in measuring the "size" of a set by the methods discussed so far, namely, by cardinality, topological density, or Lebesgue measure.


The Cantor set is a subset of the interval $[0,1]$. We will describe it by constructing its complement in $[0,1]$. The construction will proceed in stages. At the first step, let $I_{1,1}$ denote the open interval $\left(\frac{1}{3}, \frac{2}{3}\right)$. Thus, $I_{1,1}$ is the open middle third of the interval $I=[0,1]$. The second step involves performing the first step on each of the two remaining intervals of $I-I_{1,1}$. That is, we produce two open intervals, $I_{2,1}$ and $I_{2,2}$, each being the open middle third of one of the two intervals comprising $I-I_{1,1}$. At the $i^{\text {th }}$ step we produce $2^{i-1}$ open intervals, $I_{i, 1}, I_{i, 2}, \ldots, I_{i, 2^{i-1}}$, each of length $\left(\frac{1}{3}\right)^{i}$. The $(i+1)^{\text {th }}$ step consists of producing middle thirds of each of the intervals of

$$
I-\bigcup_{j=1}^{i} \bigcup_{k=1}^{2^{j-1}} I_{j, k}
$$

With $C$ denoting the Cantor set, we define its complement by

$$
I \backslash C=\bigcup_{j=1}^{\infty} \bigcup_{k=1}^{2^{j-1}} I_{j, k}
$$

Note that $C$ is a closed set and that its Lebesgue measure is 0 since

$$
\begin{aligned}
\lambda(I \backslash C) & =\frac{1}{3}+2\left(\frac{1}{3^{2}}\right)+2^{2}\left(\frac{1}{3^{3}}\right)+\cdots \\
& =\sum_{k=0}^{\infty} \frac{1}{3}\left(\frac{2}{3}\right)^{k} \\
& =1
\end{aligned}
$$

Note that $C$ is closed and $\lambda(\bar{C})=\lambda(C)=0$. Therefore $\bar{C}$ does not contain any open set since otherwise we would have $\lambda(\bar{C})>0$. This implies that $C$ is nowhere dense.

Thus, the Cantor set is small both in the sense of measure and topology. We now will determine its cardinality.

Every number $x \in[0,1]$ has a ternary expansion of the form

$$
x=\sum_{i=1}^{\infty} \frac{x_{i}}{3^{i}}
$$

where each $x_{i}$ is 0,1 , or 2 and we write $x=. x_{1} x_{2} \ldots$. This expansion is unique except when

$$
x=\frac{a}{3^{n}}
$$

where $a$ and $n$ are positive integers with $0<a<3^{n}$ and where 3 does not divide $a$. In this case $x$ has the form

$$
x=\sum_{i=1}^{n} \frac{x_{i}}{3^{i}}
$$

where $x_{i}$ is either 1 or 2 . If $x_{n}=2$, we will use this expression to represent $x$. However, if $x_{n}=1$, we will use the following representation for $x$ :

$$
x=\frac{x_{1}}{3}+\frac{x_{2}}{3^{2}}+\cdots+\frac{x_{n-1}}{3^{n-1}}+\frac{0}{3^{n}}+\sum_{i=n+1}^{\infty} \frac{2}{3^{i}}
$$

Thus, with this convention, each number $x \in[0,1]$ has a unique ternary expansion.
Let $x \in I$ and consider its ternary expansion $x=. x_{1} x_{2} \ldots$, bearing in mind the convention we have adopted above. Observe that $x \notin I_{1,1}$ if and only if $x_{1} \neq 1$. Also, if $x_{1} \neq 1$, then $x \notin I_{2,1} \cup I_{2,2}$ if and only if $x_{2} \neq 1$. Continuing in this way, we see that $x \in C$ if and only if $x_{i} \neq 1$ for each positive integer $i$. Thus, there is a one-one correspondence between elements of $C$ and all sequences $\left\{x_{i}\right\}$ where each $x_{i}$ is either 0 or 2 . The cardinality of the latter is $2^{\aleph_{0}}$ which, in view of Theorem 2.30 , is $c$.

The Cantor construction is very general and its variations lead to many interesting constructions. For example, if $0<\alpha<1$, it is possible to produce a Cantor-type set $C_{\alpha}$ in $[0,1]$ whose Lebesgue measure is $1-\alpha$. The method of
construction is the same as above, except that at the $i^{\text {th }}$ step, each of the intervals removed has length $\alpha 3^{-i}$. We leave it as an exercise to show that $C_{\alpha}$ is nowhere dense and has cardinality $c$.

## Exercises for Section 4.4

1. Prove that the Cantor-type set $C_{\alpha}$ described at the end of Section 4.4 is nowhere dense, has cardinality $c$, and has Lebesgue measure $1-\alpha$.
2. Construct an open set $U \subset[0,1]$ such that $U$ is dense in $[0,1], \lambda(U)<1$, and that $\lambda(U \cap(a, b))>0$ for any interval $(a, b) \subset[0,1]$.
3. Consider the Cantor-type set $C(\gamma)$ constructed in Section 4.8. Show that this set has the same properties as the standard Cantor set; namely, it has measure zero, it is nowhere dense, and has cardinality $c$.
4. Prove that the family of Borel subsets of $\mathbb{R}$ has cardinality $c$. From this deduce the existence of a Lebesgue measurable set which is not a Borel set.
5 . Let $E$ be the set of numbers in $[0,1]$ whose ternary expansions have only finitely many 1's. Prove that $\lambda(E)=0$.

### 4.5. Existence of Nonmeasurable Sets

The existence of a subset of $\mathbb{R}$ that is not Lebesgue measurable is intertwined with the fundamentals of set theory. Vitali showed that if the Axiom of Choice is accepted, then it is possible to establish the existence of nonmeasurable sets. However, in 1970, Solovay proved that using the usual axioms of set theory, but excluding the Axiom of Choice, it is impossible to prove the existence of a nonmeasurable set.
4.27. Theorem. There exists a set $E \subset \mathbb{R}$ that is not Lebesgue measurable.

Proof. We define a relation on elements of the real line by saying that $x$ and $y$ are equivalent (written $x \sim y$ ) if $x-y$ is a rational number. It is easily verified that $\sim$ is an equivalence relation as defined in Definition 1.1. Therefore, the real numbers are decomposed into disjoint equivalence classes. Denote the equivalence class that contains $x$ by $E_{x}$. Note that if $x$ is rational, then $E_{x}$ contains all rational numbers. Note also that each equivalence class is countable and therefore, since $\mathbb{R}$ is uncountable, there must be an uncountable number of equivalence classes. We now appeal to the Axiom of Choice, Proposition 1.4, to assert the existence of a set $S$ such that for each equivalence class $E, S \cap E$ consists precisely of one point. If $x$ and $y$ are arbitrary elements of $S$, then $x-y$ is an irrational number, for otherwise they would belong to the same equivalence class, contrary to the definition of $S$. Thus, the set of differences, defined by

$$
D_{S}:=\{x-y: x, y \in S\}
$$

is a subset of the irrational numbers and therefore cannot contain any interval. Since the Lebesgue outer measure of any set is invariant under translation and $\mathbb{R}$ is the union of the translates of $S$ by every rational number, it follows that $\lambda^{*}(S) \neq 0$. Thus, if $S$ were a measurable set, we would have $\lambda(S)>0$. If $\lambda(S)<\infty$, then Lemma 4.28 is contradicted since $D_{S}$ can not contain an interval. If $\lambda(S)=\infty$ then there exists a closed interval $I$ such that $0<\lambda(S \cap I)<\infty$ and $S \cap I$ is measurable. But this contradicts Lemma 4.28, since $D_{S \cap I} \subset D_{S}$ can not contain an interval.
4.28. Lemma. If $S \subset \mathbb{R}$ is a Lebesgue measurable set of positive and finite measure, then the set of differences $D_{S}:=\{x-y: x, y \in S\}$ contains an interval.

Proof. For each $\varepsilon>0$, there is an open set $U \supset S$ with $\lambda(U)<(1+\varepsilon) \lambda(S)$. Now $U$ is the union of a countable number of disjoint, open intervals,

$$
U=\bigcup_{k=1}^{\infty} I_{k} .
$$

Therefore,

$$
S=\bigcup_{k=1}^{\infty} S \cap I_{k} \quad \text { and } \quad \lambda(S)=\sum_{k=1}^{\infty} \lambda\left(S \cap I_{k}\right)
$$

Since $\lambda(U)=\sum_{k=1}^{\infty} \lambda\left(I_{k}\right)<(1+\varepsilon) \lambda(S)=(1+\varepsilon) \sum_{k=1}^{\infty} \lambda\left(S \cap I_{k}\right)$, it follows that $\lambda\left(I_{k_{0}}\right)<(1+\varepsilon) \lambda\left(S \cap I_{k_{0}}\right)$ for some $k_{0}$. With the choice of $\varepsilon=\frac{1}{3}$, we have

$$
\begin{equation*}
\lambda\left(S \cap I_{k_{0}}\right)>\frac{3}{4} \lambda\left(I_{k_{0}}\right) \tag{4.21}
\end{equation*}
$$

Now select any number $t$ with $0<|t|<\frac{1}{2} \lambda\left(I_{k_{0}}\right)$ and consider the translate of the set $S \cap I_{k_{0}}$ by $t$, denoted by $\left(S \cap I_{k_{0}}\right)+t$. Then $\left(S \cap I_{k_{0}}\right) \cup\left(\left(S \cap I_{k_{0}}\right)+t\right)$ is contained within an interval of length less than $\frac{3}{2} \lambda\left(I_{k_{0}}\right)$. Using the fact that the Lebesgue measure of a set remains unchanged under a translation, we conclude that the sets $S \cap I_{k_{0}}$ and $\left(S \cap I_{k_{0}}\right)+t$ must intersect, for otherwise we would contradict (4.21). This means that for each $t$ with $|t|<\frac{1}{2} \lambda\left(I_{k_{0}}\right)$, there are points $x, y \in S \cap I_{k_{0}}$ such that $x-y=t$. That is, the set

$$
D_{S} \supset\left\{x-y: x, y \in S \cap I_{k_{0}}\right\}
$$

contains an open interval centered at the origin of length $\lambda\left(I_{k_{0}}\right)$.

## Exercises for Section 4.5

1. Referring to proof of Theorem 4.27, prove that any subset of $\mathbb{R}$ with positive outer Lebesgue measure contains a nonmeasurable subset.

### 4.6. Lebesgue-Stieltjes Measure

Lebesgue-Stieltjes measure on $\mathbb{R}$ is another important outer measure that is often encountered in applications. A Lebesgue-Stieltjes measure is generated by a nondecreasing function, $f$, and its definition differs from Lebesgue measure in that the length of an interval appearing in the definition of Lebesgue measure is replaced by the oscillation of $f$ over that interval. We will show that it is a Carathéodory outer measure.

Lebesgue measure is defined by using the primitive concept of volume in $\mathbb{R}^{n}$. In $\mathbb{R}$, the length of a closed interval is used. If $f$ is a nondecreasing function defined on $\mathbb{R}$, then the "length" of a half-open interval $(a, b]$, denoted by $\alpha_{f}((a, b])$, can be defined by

$$
\begin{equation*}
\alpha_{f}((a, b])=f(b)-f(a) \tag{4.22}
\end{equation*}
$$

Based on this notion of length, a measure analogous to Lebesgue measure can be generated. This establishes an important connection between measures on $\mathbb{R}$ and monotone functions. To make this connection precise, it is necessary to use halfopen intervals in (4.22) rather than closed intervals. It is also possible to develop this procedure in $\mathbb{R}^{n}$, but it becomes more complicated, cf. [Sa].
4.29. Definition. The Lebesgue-Stieltjes outer measure of an arbitrary set $E \subset \mathbb{R}$ is defined by

$$
\begin{equation*}
\lambda_{f}^{*}(E)=\inf \left\{\sum_{h_{k} \in \mathcal{F}} \alpha_{f}\left(h_{k}\right)\right\} \tag{4.23}
\end{equation*}
$$

where the infimum is taken over all countable collections $\mathcal{F}$ of half-open intervals $h_{k}$ of the form $\left(a_{k}, b_{k}\right]$ such that

$$
E \subset \bigcup_{h_{k} \in \mathcal{F}} h_{k}
$$

Later in this section, we will show that there is an identification between LebesgueStieltjes measures and nondecreasing, right-continuous functions. This explains why we use half-open intervals of the form $(a, b]$. We could have chosen intervals of the form $[a, b)$ and then we would show that the corresponding Lebesgue-Stieltjes measure could be identified with a left-continuous function.
4.30. Remark. Also, observe that the length of each interval $\left(a_{k}, b_{k}\right]$ that appears in (4.23) can be assumed to be arbitrarily small because

$$
\alpha_{f}((a, b])=f(b)-f(a)=\sum_{k=1}^{N}\left[f\left(a_{k}\right)-f\left(a_{k-1}\right)\right]=\sum_{k=1}^{N} \alpha_{f}\left(\left(a_{k-1}, a_{k}\right]\right)
$$

whenever $a=a_{0}<a_{1}<\cdots<a_{N}=b$.
4.31. Theorem. If $f: \mathbb{R} \rightarrow \mathbb{R}$ is a nondecreasing function, then $\lambda_{f}^{*}$ is a Carathéodory outer measure on $\mathbb{R}$.

Proof. Referring to Definitions 4.1 and 4.15, we need only show that $\lambda_{f}^{*}$ is monotone, countably subadditive, and satisfies property (4.15). Verification of the remaining properties is elementary.

For the proof of monotonicity, let $A_{1} \subset A_{2}$ be arbitrary sets in $\mathbb{R}$ and assume, without loss of generality, that $\lambda_{f}^{*}\left(A_{2}\right)<\infty$. Choose $\varepsilon>0$ and consider a countable family of half-open intervals $h_{k}=\left(a_{k}, b_{k}\right]$ such that

$$
A_{2} \subset \bigcup_{k=1}^{\infty} h_{k} \text { and } \sum_{k=1}^{\infty} \alpha_{f}\left(h_{k}\right) \leq \lambda_{f}^{*}\left(A_{2}\right)+\varepsilon
$$

Then, since $A_{1} \subset \cup_{k=1}^{\infty} h_{k}$,

$$
\lambda_{f}^{*}\left(A_{1}\right) \leq \sum_{k=1}^{\infty} \alpha_{f}\left(h_{k}\right) \leq \lambda_{f}^{*}\left(A_{2}\right)+\varepsilon
$$

which establishes the desired inequality since $\varepsilon$ is arbitrary.
The proof of countable subadditivity is virtually identical to the proof of the corresponding result for Lebesgue measure given in Theorem 4.23 and thus will not be repeated here.

Similarly, the proof of property (4.9) of Definition 4.15, runs parallel to the one given in the proof of Theorem 4.23 for Lebesgue measure. Indeed, by Remark 4.30, we can may assume that the length of each $\left(a_{k}, b_{k}\right]$ is less than $\boldsymbol{d}(A, B)$.

Now that we know that $\lambda_{f}^{*}$ is a Carathéodory outer measure, it follows that the family of $\lambda_{f}^{*}$-measurable sets contains the Borel sets. As in the case of Lebesgue measure, we denote by $\lambda_{f}$ the measure obtained by restricting $\lambda_{f}^{*}$ to its family of measurable sets.

In the case of Lebesgue measure, we proved that $\lambda(I)=\operatorname{vol}(I)$ for all intervals $I \subset \mathbb{R}^{n}$. A natural question is whether the analogous property holds for $\lambda_{f}$.
4.32. Theorem. If $f: \mathbb{R} \rightarrow \mathbb{R}$ is nondecreasing and right-continuous, then

$$
\lambda_{f}((a, b])=f(b)-f(a)
$$

Proof. The proof is similar to that of Theorem 4.22 and, as in that situation, it suffices to show

$$
\lambda_{f}((a, b]) \geq f(b)-f(a)
$$

Let $\varepsilon>0$ and select a cover of $(a, b]$ by a countable family of half-open intervals, $\left(a_{i}, b_{i}\right]$ such that

$$
\begin{equation*}
\sum_{i=1}^{\infty} f\left(b_{i}\right)-f\left(a_{i}\right)<\lambda_{f}((a, b])+\varepsilon . \tag{4.24}
\end{equation*}
$$

Since $f$ is right-continuous, it follows for each $i$ that

$$
\lim _{t \rightarrow 0^{+}} \alpha_{f}\left(\left(a_{i}, b_{i}+t\right]\right)=\alpha_{f}\left(\left(a_{i}, b_{i}\right]\right)
$$

Consequently, we may replace each $\left(a_{i}, b_{i}\right.$ ] with $\left(a_{i}, b_{i}^{\prime}\right]$ where $b_{i}^{\prime}>b_{i}$ and $f\left(b_{i}^{\prime}\right)-$ $f\left(a_{i}\right)<f\left(b_{i}\right)-f\left(a_{i}\right)+\varepsilon / 2^{i}$ thus causing no essential change to (4.24), and thus allowing

$$
(a, b] \subset \bigcup_{i=1}^{\infty}\left(a_{i}, b_{i}^{\prime}\right)
$$

Let $a^{\prime} \in(a, b)$. Then

$$
\begin{equation*}
\left[a^{\prime}, b\right] \subset \bigcup_{i=1}^{\infty}\left(a_{i}, b_{i}^{\prime}\right) \tag{4.25}
\end{equation*}
$$

Let $\eta$ be the Lebesgue number of this open cover of the compact set $\left[a^{\prime}, b\right]$ (see Exercise 3.4). Partition $\left[a^{\prime}, b\right]$ into a finite number, say $m$, of intervals, each of whose length is less than $\eta$. We then have

$$
\left[a^{\prime}, b\right]=\bigcup_{k=1}^{m}\left[t_{k-1}, t_{k}\right]
$$

where $t_{0}=a^{\prime}$ and $t_{m}=b$ and each $\left[t_{k-1}, t_{k}\right]$ is contained in some element of the open cover in (4.25), say $\left(a_{i_{k}}, b_{i_{k}}^{\prime}\right]$. Furthermore, we can relabel the elements of our partition so that each $\left[t_{k-1}, t_{k}\right]$ is contained in precisely one $\left(a_{i_{k}}, b_{i_{k}}^{\prime}\right]$. Then

$$
\begin{aligned}
f(b)-f\left(a^{\prime}\right) & =\sum_{k=1}^{m} f\left(t_{k}\right)-f\left(t_{k-1}\right) \\
& \leq \sum_{k=1}^{m} f\left(b_{i_{k}}^{\prime}\right)-f\left(a_{i_{k}}\right) \\
& \leq \sum_{k=1}^{\infty} f\left(b_{i}^{\prime}\right)-f\left(a_{i}\right) \\
& \leq \lambda_{f}((a, b])+2 \varepsilon
\end{aligned}
$$

Since $\varepsilon$ is arbitrary, we have

$$
f(b)-f\left(a^{\prime}\right) \leq \lambda_{f}((a, b])
$$

Furthermore, the right continuity of $f$ implies

$$
\lim _{a^{\prime} \rightarrow a^{+}} f\left(a^{\prime}\right)=f(a)
$$

and hence

$$
f(b)-f(a) \leq \lambda_{f}((a, b])
$$

as desired.
We have just seen that a nondecreasing function $f$ gives rise to a Borel outer measure on $\mathbb{R}$. The converse is readily seen to hold, for if $\mu$ is a finite Borel outer measure on $\mathbb{R}$ (see Definition 4.14), let

$$
f(x)=\mu((-\infty, x])
$$

Then, $f$ is nondecreasing, right-continuous (see Exercise 4.1) and

$$
\mu((a, b])=f(b)-f(a) \quad \text { whenever } \quad a<b
$$

(Incidentally, this now shows why half-open intervals are used in the development.) With $f$ defined this way, note from our previous result, Theorem 4.32, that the corresponding Lebesgue- Stieltjes measure, $\lambda_{f}$, satisfies

$$
\lambda_{f}((a, b])=f(b)-f(a)
$$

thus proving that $\mu$ and $\lambda_{f}$ agree on all half-open intervals. Since every open set in $\mathbb{R}$ is a countable union of disjoint half-open intervals, it follows that $\mu$ and $\lambda_{f}$ agree on all open sets. Consequently, it seems plausible that these measures should agree on all Borel sets. In fact, this is true because both $\mu$ and $\lambda_{f}^{*}$ are outer measures with the approximation property described in Theorem 4.52 below. Consequently, we have the following result.
4.33. Theorem. Suppose $\mu$ is a finite Borel outer measure on $\mathbb{R}$ and let

$$
f(x)=\mu((-\infty, x])
$$

Then the Lebesgue- Stieltjes measure, $\lambda_{f}$, agrees with $\mu$ on all Borel sets.

## Exercises for Section 4.6

1. Suppose $\mu$ is a finite Borel measure defined on $\mathbb{R}$.

Let $f(x)=\mu((-\infty, x])$. Prove that $f$ is right continuous.
2. Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a nondecreasing function and let $\lambda_{f}$ be the Lebesgue- Stieltjes measure generated by $f$. Prove that $\lambda_{f}\left(\left\{x_{0}\right\}\right)=0$ if and only if $f$ is leftcontinuous at $x_{0}$.
3. Let $f$ be a nondecreasing function defined on $\mathbb{R}$. Define a Lebesgue-Stieltjes-type measure as follows: For $A \subset \mathbb{R}$ an arbitrary set,

$$
\begin{equation*}
\Lambda_{f}^{*}(A)=\inf \left\{\sum_{h_{k} \in \mathcal{F}}\left[f\left(b_{k}\right)-f\left(a_{k}\right)\right]\right\} \tag{4.26}
\end{equation*}
$$

where the infimum is taken over all countable collections $\mathcal{F}$ of closed intervals of the form $h_{k}:=\left[a_{k}, b_{k}\right]$ such that

$$
E \subset \bigcup_{h_{k} \in \mathcal{F}} h_{k}
$$

In other words, the definition of $\Lambda_{f}^{*}(A)$ is the same as $\lambda_{f}^{*}(A)$ except that closed intervals $\left[a_{k}, b_{k}\right]$ are used instead of half-open intervals $\left(a_{k}, b_{k}\right]$.

As in the case of Lebesgue- Stieltjes measure it can be easily seen that $\Lambda_{f}^{*}$ is a Carathédory measure. (You need not prove this).
(a) Prove that $\Lambda_{f}^{*}(A) \leq \lambda_{f}^{*}(A)$ for all sets $A \subset \mathbb{R}^{n}$.
(b) Prove that $\Lambda_{f}^{*}(B)=\lambda_{f}^{*}(B)$ for all Borel sets $B$ if $f$ is left-ontinuous.

### 4.7. Hausdorff Measure


#### Abstract

As a final illustration of a Carathéodory measure, we introduce $s$ dimensional Hausdorff (outer) measure in $\mathbb{R}^{n}$ where $s$ is any nonnegative real number. It will be shown that the only significant values of $s$ are those for which $0 \leq s \leq n$ and that for $s$ in this range, Hausdorff measure provides meaningful measurements of small sets. For example, sets of Lebesgue measure zero may have positive Hausdorff measure.


4.34. Definitions. Hausdorff measure is defined in terms of an auxiliary set function that we introduce first. Let $0 \leq s<\infty, 0<\varepsilon \leq \infty$ and let $A \subset \mathbb{R}^{n}$. Define

$$
\begin{equation*}
H_{\varepsilon}^{s}(A)=\inf \left\{\sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s}: A \subset \bigcup_{i=1}^{\infty} E_{i}, \operatorname{diam} E_{i}<\varepsilon\right\} \tag{4.27}
\end{equation*}
$$

where $\alpha(s)$ is a normalization constant defined by

$$
\alpha(s)=\frac{\pi^{\frac{s}{2}}}{\Gamma\left(\frac{s}{2}+1\right)},
$$

with

$$
\Gamma(t)=\int_{0}^{\infty} e^{-x} x^{t-1} d x, \quad 0<t<\infty
$$

It follows from definition that if $\varepsilon_{1}<\varepsilon_{2}$, then $H_{\varepsilon_{1}}^{s}(E) \geq H_{\varepsilon_{2}}^{s}(E)$. This allows the following, which is the definition of $s$-dimensional Hausdorff measure:

$$
H^{s}(A)=\lim _{\varepsilon \rightarrow 0} H_{\varepsilon}^{s}(A)=\sup _{\varepsilon>0} H_{\varepsilon}^{s}(A) .
$$

When $s$ is a positive integer, it turns out that $\alpha(s)$ is the Lebesgue measure of the unit ball in $\mathbb{R}^{s}$. This makes it possible to prove that $H^{s}$ assigns to elementary sets the value one would expect. For example, consider $n=3$. In this case $\alpha(3)=\frac{\pi^{3 / 2}}{\Gamma\left(\frac{3}{2}+1\right)}=\frac{\pi^{3 / 2}}{\Gamma\left(\frac{5}{2}\right)}=\frac{\pi^{3 / 2}}{\frac{3}{4} \pi^{1 / 2}}=\frac{4}{3} \pi$. Note that $\alpha(3) 2^{-3}(\operatorname{diam} B(x, r))^{3}=\frac{4}{3} \pi r^{3}=$ $\lambda(B(x, r))$. In fact, it can be shown that $H^{n}(B(x, r))=\lambda(B(x, r))$ for any ball $B(x, r)$ (see Exercise 1, Section 4.7). In Definition 4.27 we have fixed $n>0$ and
we have defined, for any $0 \leq s<\infty$, the $s$-dimensional Hausdorff measures $H^{s}$. However, for any $A \subset \mathbb{R}^{n}$ and $s>n$ we have $H^{s}(A)=0$. That is, $H^{s} \equiv 0$ on $\mathbb{R}^{n}$ for all $s>n$ (see Exercise 3, Section 4.7).

Before deriving the basic properties of $H^{s}$, a few observations are in order.

### 4.35. REmARK.

(i) Hausdorff measure could be defined in any metric space since the essential part of the definition depends on only the notion of diameter of a set.
(ii) The sets $E_{i}$ in the definition of $H_{\varepsilon}^{s}(A)$ are arbitrary subsets of $\mathbb{R}^{n}$. However, they could be taken to be closed sets since diam $E_{i}=\operatorname{diam} \overline{E_{i}}$.
(iii) The reason for the restriction of coverings by sets of small diameter is to produce an accurate measurement of sets that are geometrically complicated. For example, consider the set $A=\{(x, \sin (1 / x)): 0<x \leq 1\}$ in $\mathbb{R}^{2}$. We will see in Section 7.8 that $H^{1}(A)$ is the length of the set $A$, so that in this case $H^{1}(A)=\infty$ (it is an instructive exercise to show this directly from the Definition). If no restriction on the diameter of the covering sets were imposed, the measure of $A$ would be finite.
(iv) Often Hausdorff measure is defined without the inclusion of the constant $\alpha(s) 2^{-s}$. Then the resulting measure differs from our definition by a constant factor, which is not important unless one is interested in the precise value of Hausdorff measure

We now proceed to derive some of the basic properties of Hausdorff measure.
4.36. Theorem. For each nonnegative number $s, H^{s}$ is a Carathéodory outer measure.

Proof. We must show that the four conditions of Definition 4.1 are satisfied as well as condition (4.15). The first three conditions of Definition 4.1 are immediate, and so we proceed to show that $H^{s}$ is countably subadditive. For this, suppose $\left\{A_{i}\right\}$ is a sequence of sets in $\mathbb{R}^{n}$ and select sets $\left\{E_{i, j}\right\}$ such that

$$
A_{i} \subset \bigcup_{j=1}^{\infty} E_{i, j}, \quad \operatorname{diam} E_{i, j} \leq \varepsilon, \quad \sum_{j=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i, j}\right)^{s}<H_{\varepsilon}^{s}\left(A_{i}\right)+\frac{\varepsilon}{2^{i}}
$$

Then, as $i$ and $j$ range through the positive integers, the sets $\left\{E_{i, j}\right\}$ produce a countable covering of $A$ and therefore,

$$
\begin{aligned}
H_{\varepsilon}^{s}\left(\bigcup_{i=1}^{\infty} A_{i}\right) & \leq \sum_{i=1}^{\infty} \sum_{j=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i, j}\right)^{s} \\
& =\sum_{i=1}^{\infty}\left[H_{\varepsilon}^{s}\left(A_{i}\right)+\frac{\varepsilon}{2^{i}}\right]
\end{aligned}
$$

Now $H_{\varepsilon}^{s}\left(A_{i}\right) \leq H^{s}\left(A_{i}\right)$ for each $i$ so that

$$
H_{\varepsilon}^{s}\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \sum_{i=1}^{\infty} H^{s}\left(A_{i}\right)+\varepsilon
$$

Now taking limits as $\varepsilon \rightarrow 0$ we obtain

$$
H^{s}\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \sum_{i=1}^{\infty} H^{s}\left(A_{i}\right)
$$

which establishes countable subadditivity.
Now we will show that condition (4.15) is satisfied. Choose $A, B \subset \mathbb{R}^{n}$ with $\boldsymbol{d}(A, B)>0$ and let $\varepsilon$ be any positive number less than $\boldsymbol{d}(A, B)$. Let $\left\{E_{i}\right\}$ be a covering of $A \cup B$ with $\operatorname{diam} E_{i} \leq \varepsilon$. Thus no set $E_{i}$ intersects both $A$ and $B$. Let $\mathcal{A}$ be the collection of those $E_{i}$ that intersect $A$, and $\mathcal{B}$ those that intersect $B$. Then

$$
\begin{aligned}
\sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s} \geq & \sum_{E \in \mathcal{A}} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s} \\
& +\sum_{E \in \mathcal{B}} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s} \\
\geq & H_{\varepsilon}^{s}(A)+H_{\varepsilon}^{s}(B)
\end{aligned}
$$

Taking the infimum over all such coverings $\left\{E_{i}\right\}$, we obtain

$$
H_{\varepsilon}^{s}(A \cup B) \geq H_{\varepsilon}^{s}(A)+H_{\varepsilon}^{s}(B)
$$

where $\varepsilon$ is any number less than $\boldsymbol{d}(A, B)$. Finally, taking the limit as $\varepsilon \rightarrow 0$, we have

$$
H^{s}(A \cup B) \geq H^{s}(A)+H^{s}(B)
$$

Since we already established (countable) subadditivity of $H^{s}$, property (4.15) is thus established and the proof is concluded.

Since $H^{s}$ is a Carathéodory outer measure, it follows from Theorem 4.18 that all Borel sets are $H^{s}$-measurable. We next show that $H^{s}$ is, in fact, a Borel regular outer measure in the sense of Definition 4.14.
4.37. Theorem. For each $A \subset \mathbb{R}^{n}$, there exists a Borel set $B \supset A$ such that

$$
H^{s}(B)=H^{s}(A)
$$

Proof. From the previous comment, we already know that $H^{s}$ is a Borel outer measure. To show that it is a Borel regular outer measure, recall from (ii) in Remark 4.35 above that the sets $\left\{E_{i}\right\}$ in the definition of Hausdorff measure can be taken as closed sets. Suppose $A \subset \mathbb{R}^{n}$ with $H^{s}(A)<\infty$, thus implying that $H_{\varepsilon}^{s}(A)<\infty$ for all $\varepsilon>0$. Let $\left\{\varepsilon_{j}\right\}$ be a sequence of positive numbers such that $\varepsilon_{j} \rightarrow 0$, and for each positive integer $j$, choose closed sets $\left\{E_{i, j}\right\}$ such that $\operatorname{diam} E_{i, j} \leq \varepsilon_{j}, A \subset \cup_{i=1}^{\infty} E_{i, j}$, and

$$
\sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i, j}\right)^{s} \leq H_{\varepsilon_{j}}^{s}(A)+\varepsilon_{j}
$$

Set

$$
A_{j}=\bigcup_{i=1}^{\infty} E_{i, j} \quad \text { and } \quad B=\bigcap_{j=1}^{\infty} A_{j}
$$

Then $B$ is a Borel set and since $A \subset A_{j}$ for each $j$, we have $A \subset B$. Furthermore, since

$$
B \subset \bigcup_{i=1}^{\infty} E_{i, j}
$$

for each $j$, we have

$$
H_{\varepsilon_{j}}^{s}(B) \leq \sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i, j}\right)^{s} \leq H_{\varepsilon_{j}}^{s}(A)+\varepsilon_{j}
$$

Since $\varepsilon_{j} \rightarrow 0$ as $j \rightarrow \infty$, we obtain $H^{s}(B) \leq H^{s}(A)$. But $A \subset B$ so that we have $H^{s}(A)=H^{s}(B)$.
4.38. Remark. The preceding result can be improved. In fact, there is a $G_{\delta}$ set $G$ containing $A$ such that $H^{s}(G)=H^{s}(A)$; see Exercise 8, Section 4.7.
4.39. Theorem. Suppose $A \subset \mathbb{R}^{n}$ and $0 \leq s<t<\infty$. Then
(i) If $H^{s}(A)<\infty$ then $H^{t}(A)=0$
(ii) If $H^{t}(A)>0$ then $H^{s}(A)=\infty$

Proof. We need only prove (i) because (ii) is simply a restatement of (i). We state (ii) only to emphasize its importance.

For the proof of (i), choose $\varepsilon>0$ and a covering of $A$ by sets $\left\{E_{i}\right\}$ with $\operatorname{diam} E_{i}<\varepsilon$ such that

$$
\sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s} \leq H_{\varepsilon}^{s}(A)+1 \leq H^{s}(A)+1
$$

Then

$$
\begin{aligned}
H_{\varepsilon}^{t}(A) & \leq \sum_{i=1}^{\infty} \alpha(t) 2^{-t}\left(\operatorname{diam} E_{i}\right)^{t} \\
& =\frac{\alpha(t)}{\alpha(s)} 2^{s-t} \sum_{i=1}^{\infty} \alpha(s) 2^{-s}\left(\operatorname{diam} E_{i}\right)^{s}\left(\operatorname{diam} E_{i}\right)^{t-s} \\
& \leq \frac{\alpha(t)}{\alpha(s)} 2^{s-t} \varepsilon^{t-s}\left[H^{s}(A)+1\right]
\end{aligned}
$$

Now let $\varepsilon \rightarrow 0$ to obtain $H^{t}(A)=0$.
4.40. Definition. The Hausdorff dimension of an arbitrary set $A \subset \mathbb{R}^{n}$ is that number $0 \leq \delta_{A} \leq n$ such that

$$
\delta_{A}=\inf \left\{t: H^{t}(A)=0\right\}=\sup \left\{s: H^{s}(A)=\infty\right\}
$$

In other words, the Hausdorff dimension $\delta_{A}$ is that unique number such that

$$
\begin{aligned}
& s<\delta_{A} \text { implies } H^{s}(A)=\infty \\
& t>\delta_{A} \text { implies } H^{t}(A)=0
\end{aligned}
$$

The existence and uniqueness of $\delta_{A}$ follows directly from Theorem 4.39.
4.41. Remark. If $s=\delta_{A}$ then one of the following three possibilities has to occur: $H^{s}(A)=0, H^{s}(A)=\infty$ or $0<H^{s}(A)<\infty$. On the other hand, if $0<H^{s}(A)<\infty$, then it follows from Theorem 4.39 that $\delta_{A}=s$.

The notion of Hausdorff dimension is not very intuitive. Indeed, the Hausdorff dimension of a set need not be an integer. Moreover, if the dimension of a set is an integer $k$, the set need not resemble a " $k$-dimensional surface" in any usual sense. See Falconer [26] or Federer [27] for examples of pathological Cantor-like sets with integer Hausdorff dimension. However, we can at least be reassured by the fact that the Hausdorff dimension of an open set $U \subset \mathbb{R}^{n}$ is $n$. To verify this, it is sufficient to assume that $U$ is bounded and to prove that

$$
\begin{equation*}
0<H^{n}(U)<\infty \tag{4.28}
\end{equation*}
$$

Exercise 2, Section 4.7, deals with the proof of this. Also, it is clear that any countable set has Hausdorff dimension zero; however, there are uncountable sets with dimension zero (see Exercise 7, Section 4.7).

## Exercises for Section 4.7

1. If $A \subset \mathbb{R}$ is an arbitrary set, show that $H^{1}(A)=\lambda^{*}(A)$.
4.42. Remark. In this problem, you will see the importance of the constant that appears in the definition of Hausdorff measure. The constant $\alpha(s)$ that appears in the definition of $H^{s}$-measure is equal to 2 when $s=1$. That is, $\alpha(1)=2$, and therefore, the definition of $H^{1}(A)$ can be written as

$$
H^{1}(A)=\lim _{\varepsilon \rightarrow 0} H_{\varepsilon}^{1}(A)
$$

where

$$
H_{\varepsilon}^{1}(A)=\inf \left\{\sum_{i=1}^{\infty} \operatorname{diam} E_{i}: A \subset \bigcup_{i=1}^{\infty} E_{i} \subset \mathbb{R}, \operatorname{diam} E_{i}<\varepsilon\right\}
$$

This result is also true in $\mathbb{R}^{n}$ but more difficult to prove. The isodiametric inequality $\lambda^{*}(A) \leq \alpha(n)\left(\frac{\operatorname{diam} A}{2}\right)^{n}$ (whose proof is omitted in this book), can be used to prove that $H^{n}(A)=\lambda^{*}(A)$ for all $A \subset \mathbb{R}^{n}$.
2. For $A \subset \mathbb{R}^{n}$, use the isodiametric inequality introduced in exercise 4.1 to show that

$$
\lambda^{*}(A) \leq H^{n}(A) \leq \alpha(n)\left(\frac{\sqrt{n}}{2}\right)^{n} \lambda^{*}(A)
$$

3. Show that $H^{s} \equiv 0$ on $\mathbb{R}^{n}$ for all $s>n$.
4. Let $C \subset \mathbb{R}^{2}$ denote the circle of radius 1 and consider $C$ as a topological space with the topology induced from $\mathbb{R}^{2}$. Define an outer measure $\mathcal{H}$ on $C$ by

$$
\mathcal{H}(A):=\frac{1}{2 \pi} H^{1}(A) \quad \text { for any set } A \subset C
$$

Later in the course, we will prove that $H^{1}(C)=2 \pi$. Thus, you may assume that. Show that
(i) $\mathcal{H}(C)=1$.
(ii) Prove that $\mathcal{H}$ is a Borel regular outer measure.
(iii) Prove that $\mathcal{H}$ is rotationally invariant; that is, prove for any set $A \subset C$ that $\mathcal{H}(A)=\mathcal{H}\left(A^{\prime}\right)$ where $A^{\prime}$ is obtained by rotating $A$ through an arbitrary angle.
(iv) Prove that $\mathcal{H}$ is the only outer measure defined on $C$ that satisfies the previous 3 conditions.
5. Another Hausdorff-type measure that is frequently used is Hausdorff spherical measure, $H_{S}^{s}$. It is defined in the same way as Hausdorff measure (Definition 4.34) except that the sets $E_{i}$ are replaced by $n$-balls. Clearly, $H^{s}(E) \leq H_{S}^{s}(E)$ for any set $E \subset \mathbb{R}^{n}$. Prove that $H_{S}^{s}(E) \leq 2^{s} H^{s}(E)$ for any set $E \subset \mathbb{R}^{n}$.
6. Prove that a countable set $A \subset \mathbb{R}^{n}$ has Hausdorff dimension 0 . The following problem shows that the converse is not tue.
7. Let $S=\left\{a_{i}\right\}$ be any sequence of real numbers in $(0,1 / 2)$. We now will construct a Cantor set $C(S)$ similar in construction to that of $C(\lambda)$ except the length of the intervals $I_{k, j}$ at the $k$ th stage will not be a constant multiple of those in the preceding stage. Instead, we proceed as follows: Define $I_{0,1}=[0,1]$ and then define the both intervals $I_{1,1}, I_{1,2}$ to have length $a_{1}$. Proceeding inductively, the intervals at the $k^{\text {th }}, I_{k, i}$, will have length $a_{k} l\left(I_{k-1, i}\right)$. Consequently, at the $k^{\text {th }}$ stage, we obtain $2^{k}$ intervals $I_{k, j}$ each of length

$$
s_{k}=a_{1} a_{2} \cdots a_{k}
$$

It can be easily verified that the resulting Cantor set $C(S)$ has cardinality $c$ and is nowhere dense.

The focus of this problem is to determine the Hausdorff dimension of $C(S)$. For this purpose, consider the function defined on $(0, \infty)$

$$
h(r):= \begin{cases}\frac{r^{s}}{\log (1 / r)} & \text { when } 0 \leq s<1  \tag{4.29}\\ r^{s} \log (1 / r) & \text { where } 0<s \leq 1\end{cases}
$$

Note that $h$ is increasing and $\lim _{r \rightarrow 0} h(r)=0$. Corresponding to this function, we will construct a Cantor set $C\left(S_{h}\right)$ that will have interesting properties. We will select inductively numbers $a_{1}, a_{2}, \ldots$ such that

$$
\begin{equation*}
h\left(s_{k}\right)=2^{-k} \tag{4.30}
\end{equation*}
$$

That is, $a_{1}$ is chosen so that $h\left(a_{1}\right)=1 / 2$, i.e., $a_{1}=h^{-1}(1 / 2)$. Now that $a_{1}$ has been chosen, let $a_{2}$ be that number such that $h\left(a_{1} a_{2}\right)=1 / 2^{2}$. In this way, we can choose a sequence $S_{h}:=\left\{a_{1}, a_{2}, \ldots\right\}$ such that (4.30) is satisfied. Now consider the following Hausdorff-type measure:

$$
H_{\varepsilon}^{h}(A):=\inf \left\{\sum_{i=1}^{\infty} h\left(\operatorname{diam} E_{i}\right): A \subset \bigcup_{i=1}^{\infty} E_{i} \subset \mathbb{R}^{n}, h\left(\operatorname{diam} E_{i}\right)<\varepsilon\right\}
$$

and

$$
H^{h}(A):=\lim _{\varepsilon \rightarrow 0} H_{\varepsilon}^{h}(A)
$$

With the Cantor set $C\left(S_{h}\right)$ that was constructed above, it follows that

$$
\begin{equation*}
\frac{1}{4} \leq H^{h}\left(C\left(S_{h}\right)\right) \leq 1 \tag{4.31}
\end{equation*}
$$

The proof of this proceeds precisely the same way as in Section 4.8.
(a) With $s=0$, our function $h$ in (4.29) becomes $h(r)=1 / \log (1 / r)$ and we obtain a corresponding Cantor set $C\left(S_{h}\right)$. With the help of (4.31) prove that the Hausdorff dimension of $C\left(S_{h}\right)$ is zero, thus showing that the converse of Problem 1 is not true.
(b) Now take $s=1$ and then our function $h$ in (4.29) becomes $h(r)=r \log (1 / r)$ and again we obtain a corresponding Cantor set $C\left(S_{h}\right)$. Prove that the Hausdorff dimension of $C\left(S_{h}\right)$ is 1 which shows that there are sets other than intervals in $\mathbb{R}$ that have dimension 1.
8. For each arbitrary set $A \subset \mathbb{R}^{n}$, prove that there exists a $G_{\delta}$ set $B \supset A$ such that

$$
\begin{equation*}
H^{s}(B)=H^{s}(A) \tag{4.32}
\end{equation*}
$$

4.43. Remark. This result shows that all three of our primary measures, namely, Lebesgue measure, Lebesgue- Stieltjes measure and Hausdorff measure, share the same important regularity property (4.32).
9. If $A \subset \mathbb{R}^{n}$ is an arbitrary set and $0 \leq t \leq n$, prove that if $H_{\varepsilon}^{t}(A)=0$ for some $0<\varepsilon \leq \infty$, then $H^{t}(A)=0$.
10. Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be Lipschitz (see Definition 3.65), $E \subset \mathbb{R}^{n}, 0 \leq s<\infty$. Then

$$
\mathcal{H}^{s}(f(A)) \leq C_{f}^{s} \mathcal{H}^{s}(A)
$$

### 4.8. Hausdorff Dimension of Cantor Sets

In this section the Hausdorff dimension of Cantor sets will be determined. Note that for $H^{1}$ defined in $\mathbb{R}$ that the constant $\alpha(s) 2^{-s}$ in (4.34) equals 1.
4.44. Definition. [General Cantor Set] Let $0<\gamma<1 / 2$ and denote $I_{0,1}=$ $[0,1]$. Let $I_{1,1}$ and $I_{1,2}$ denote the intervals $[0, \gamma]$ and $[1-\gamma, 1]$ respectively. They result by deleting the open middle interval of length $1-2 \gamma$. At the next stage, delete the open middle interval of length $\gamma(1-2 \gamma)$ of each of the intervals $I_{1,1}$ and $I_{1,2}$. There remains $2^{2}$ closed intervals each of length $\gamma^{2}$. Continuing this process, at the $k^{t h}$ stage there are $2^{k}$ closed intervals each of length $\gamma^{k}$. Denote these intervals by $I_{k, 1}, \ldots, I_{k, 2^{k}}$. We define the generalized Cantor set as

$$
C(\gamma)=\bigcap_{k=0}^{\infty} \bigcup_{j=1}^{2^{k}} I_{k, j}
$$

Note that $C(1 / 3)$ is the Cantor set discussed in Section 4.4.

$$
I_{0,1}
$$

$\qquad$
$\qquad$

| $I_{2,1}$ | $\underline{I_{2,2}}$ |
| :--- | :--- |



Since $C(\gamma) \subset \bigcup_{j=1}^{2^{k}} I_{k, j}$ for each $k$, it follows that

$$
H_{\gamma^{k}}^{s}(C(\gamma)) \leq \sum_{j=1}^{2^{k}} l\left(I_{k, j}\right)^{s}=2^{k} \gamma^{k s}=\left(2 \gamma^{s}\right)^{k},
$$

where

$$
l\left(I_{k, j}\right) \text { denotes the length of } I_{k, j} .
$$

If $s$ is chosen so that $2 \gamma^{s}=1$, (if $s=\log 2 / \log (1 / \gamma)$ ) we have

$$
\begin{equation*}
H^{s}(C(\gamma))=\lim _{k \rightarrow \infty} H_{\gamma^{k}}^{s}(C(\gamma)) \leq 1 \tag{4.33}
\end{equation*}
$$

It is important to observe that our choice of $s$ implies that the sum of the $s$-power of the lengths of the intervals at any stage is one; that is,

$$
\begin{equation*}
\sum_{j=1}^{2^{k}} l\left(I_{k, j}\right)^{s}=1 . \tag{4.34}
\end{equation*}
$$

Next, we show that $H^{s}(C(\gamma)) \geq 1 / 4$ which, along with (4.33), implies that the Hausdorff dimension of $C(\gamma)=\log (2) / \log (1 / \gamma)$. We will establish this by showing that if

$$
C(\gamma) \subset \bigcup_{i=1}^{\infty} J_{i}
$$

is an open covering $C(\gamma)$ by intervals $J_{i}$ then

$$
\begin{equation*}
\sum_{i=1}^{\infty} l\left(J_{i}\right)^{s} \geq \frac{1}{4} \tag{4.35}
\end{equation*}
$$

Since this is an open cover of the compact set $C(\gamma)$, we can employ the Lebesgue number of this covering to conclude that each interval $I_{k, j}$ of the $k^{\text {th }}$ stage is contained in some $J_{i}$ provided $k$ is sufficiently large. We will show for any open interval $I$ and any fixed $\ell$, that

$$
\begin{equation*}
\sum_{I_{\ell, i} \subset I} l\left(I_{\ell, i}\right)^{s} \leq 4 l(I)^{s} \tag{4.36}
\end{equation*}
$$

This will establish (4.35) since

$$
\begin{aligned}
4 \sum_{i} l\left(J_{i}\right)^{s} & \geq \sum_{i} \sum_{I_{k, j} \subset J_{i}} l\left(I_{k, j}\right)^{s} \quad \text { by (4.36) } \\
& \geq \sum_{j=1}^{\infty} l\left(I_{k, j}\right)^{s}=1 \quad \text { because } \bigcup_{i=1}^{2^{k}} I_{k, i} \subset \bigcup_{i=1}^{\infty} J_{i} \text { and by }
\end{aligned}
$$

To verify (4.36), assume that $I$ contains some interval $I_{\ell, i}$ from the $\ell^{\text {th }}$ stage. and let $k$ denote the smallest integer for which $I$ contains some interval $I_{k, j}$ from the $k^{t h}$ stage. Then $k \leq \ell$. By considering the construction of our set $C(\gamma)$, it follows that no more than 4 intervals from the $k^{t h}$ stage can intersect $I$ for otherwise, $I$ would contain some $I_{k-1, i}$. Call the intervals $I_{k, k_{m}}, m=1,2,3,4$. Thus,

$$
\begin{aligned}
4 l(I)^{s} \geq \sum_{m=1}^{4} l\left(I_{k, k_{m}}\right)^{s} & =\sum_{m=1}^{4} \sum_{I_{\ell, i} \subset I_{k, k_{m}}} l\left(I_{\ell, i}\right)^{s} \\
& \geq \sum_{I_{\ell, i} \subset I} l\left(I_{\ell, i}\right)^{s} .
\end{aligned}
$$

which establishes (4.36). This proves that the dimension of $C(\gamma)=\log 2 / \log (1 / \gamma)$.
4.45. Remark. It can be shown that (4.35) can be improved to read

$$
\begin{equation*}
\sum_{i}^{\infty} l\left(J_{i}\right)^{s} \geq 1 \tag{4.37}
\end{equation*}
$$

which implies the precise result $H^{s}(C(\gamma))=1$ if

$$
s=\frac{\log 2}{\log (1 / \gamma)}
$$

4.46. Remark. The Cantor sets $C(\gamma)$ are prototypical examples of sets that possess self-similar properties. A set is self-similar if it can be decomposed into parts which are geometrically similar to the whole set. For example, the sets $C(\gamma) \cap$ $[0, \gamma]$ and $C(\gamma) \cap[1-\gamma, 1]$ when magnified by the factor $1 / \gamma$ yield a translate of $C(\gamma)$. Self-similarity is the characteristic property of fractals.

### 4.9. Measures on Abstract Spaces

Given an arbitrary set $X$ and a $\sigma$-algebra, $\mathcal{M}$, of subsets of $X$, a nonnegative, countably additive set function defined on $\mathcal{M}$ is called a measure. In this section we extract the properties of outer measures when restricted to their measurable sets.

Before proceeding, recall the development of the first three sections of this chapter. We began with the concept of an outer measure on an arbitrary set $X$ and proved that the family of measurable sets forms a $\sigma$-algebra. Furthermore, we showed that the outer measure is countably additive on measurable sets. In order to ensure that there are situations in which the family of measurable sets is large, we investigated Carathéodory outer measures on a metric space and established that their measurable sets always contain the Borel sets. We then introduced Lebesgue measure as the primary example of a Carathéodory outer measure. In this development, we begin to see that countable additivity plays a central and indispensable role and thus, we now call upon a common practice in mathematics of placing a crucial concept in an abstract setting in order to isolate it from the clutter and distractions of extraneous ideas. We begin with the following definition:
4.47. Definition. Let $X$ be a set and $\mathcal{M}$ a $\sigma$-algebra of subsets of $X$. A measure on $\mathcal{M}$ is a function $\mu: \mathcal{M} \rightarrow[0, \infty]$ satisfying the properties
(i) $\mu(\emptyset)=0$,
(ii) if $\left\{E_{i}\right\}$ is a sequence of disjoint sets in $\mathcal{M}$, then

$$
\mu\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\sum_{i=1}^{\infty} \mu\left(E_{i}\right)
$$

Thus, a measure is a countably additive set function defined on $\mathcal{M}$. Sometimes the notion of finite additivity is useful. It states that
$(\text { ii })^{\prime}$ If $E_{1}, E_{2}, \ldots, E_{k}$ is any finite family of disjoint sets in $\mathcal{M}$, then

$$
\mu\left(\bigcup_{i=1}^{k} E_{i}\right)=\sum_{i=1}^{k} \mu\left(E_{i}\right) .
$$

If $\mu$ satisfies (i) and (ii'), but not necessarily (ii), $\mu$ is called a finitely additive measure. The triple $(X, \mathcal{M}, \mu)$ is called a measure space and the sets that constitute $\mathcal{M}$ are called measurable sets. To be precise these sets should be referred to as $\mathcal{M}$-measurable, to indicate their dependence on $\mathcal{M}$. However, in most situations, it will be clear from the context which $\sigma$-algebra is intended and thus, the more involved notation will not be required. In case $\mathcal{M}$ constitutes the family of Borel sets in a metric space $X, \mu$ is called a Borel measure. A measure $\mu$ is said to be finite if $\mu(X)<\infty$ and $\sigma$-finite if $X$ can be written as $X=\cup_{i=1}^{\infty} E_{i}$
where $\mu\left(E_{i}\right)<\infty$ for each $i$. A measure $\mu$ with the property that all subsets of sets of $\mu$-measure zero are measurable, is said to be complete and $(X, \mathcal{M}, \mu)$ is called a complete measure space. A Borel measure on a topological space $X$ that is finite on compact sets is called a Radon measure. Thus, Lebesgue measure on $\mathbb{R}^{n}$ is a Radon measure, but $s$-dimensional Hausdorff measure, $0 \leq s<n$, is not (see Theorem 4.63 for regularity properties of Borel measures).

We emphasize that the notation $\mu(E)$ implies that $E$ is an element of $\mathcal{M}$, since $\mu$ is defined only on $\mathcal{M}$. Thus, when we write $\mu\left(E_{i}\right)$ as in the definition above, it should be understood that the sets $E_{i}$ are necessarily elements of $\mathcal{M}$.
4.48. Examples. Here are some examples of measures.
(i) $\left(\mathbb{R}^{n}, \mathcal{M}, \lambda\right)$ where $\lambda$ is Lebesgue measure and $\mathcal{M}$ is the family of Lebesgue measurable sets.
(ii) $(X, \mathcal{M}, \varphi)$ where $\varphi$ is an outer measure on an abstract set $X$ and $\mathcal{M}$ is the family of $\varphi$-measurable sets.
(iii) $\left(X, \mathcal{M}, \delta_{x_{0}}\right)$ where $X$ is an arbitrary set and $\delta_{x_{0}}$ is an outer measure defined by

$$
\delta_{x_{0}}(E)= \begin{cases}1 & \text { if } x_{0} \in E \\ 0 & \text { if } x_{0} \notin E\end{cases}
$$

The point $x_{0} \in X$ is selected arbitrarily. It can easily be shown that all subsets of $X$ are $\delta_{x_{0}}$-measurable and therefore, $\mathcal{M}$ is taken as the family of all subsets of $X$.
(iv) $(\mathbb{R}, \mathcal{M}, \mu)$ where $\mathcal{M}$ is the family of all Lebesgue measurable sets, $x_{0} \in \mathbb{R}$ and $\mu$ is defined by

$$
\mu(E)=\lambda\left(E \backslash\left\{x_{0}\right\}\right)+\delta_{x_{0}}(E)
$$

whenever $E \in \mathcal{M}$.
(v) $(X, \mathcal{M}, \mu)$ where $\mathcal{M}$ is the family of all subsets of an arbitrary space $X$ and where $\mu(E)$ is defined as the number (possibly infinite) of points in $E \in \mathcal{M}$.

The proof of Corollary 4.12 utilized only those properties of an outer measure that an abstract measure possesses and therefore, most of the following do not require a proof.
4.49. Theorem. Let $(X, \mathcal{M}, \mu)$ be a measure space and suppose $\left\{E_{i}\right\}$ is a sequence of sets in $\mathcal{M}$.
(i) (Monotonicity) If $E_{1} \subset E_{2}$, then $\mu\left(E_{1}\right) \leq \mu\left(E_{2}\right)$.
(ii) (Subtractivity) If $E_{1} \subset E_{2}$ and $\mu\left(E_{1}\right)<\infty$, then $\mu\left(E_{2}-E_{1}\right)=\mu\left(E_{2}\right)-\mu\left(E_{1}\right)$.
(iii) (Countable Subadditivity)

$$
\mu\left(\bigcup_{i=1}^{\infty} E_{i}\right) \leq \sum_{i=1}^{\infty} \mu\left(E_{i}\right)
$$

(iv) (Continuity from the left) If $\left\{E_{i}\right\}$ is an increasing sequence of sets, that is, if $E_{i} \subset E_{i+1}$ for each $i$, then

$$
\mu\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\mu\left(\lim _{i \rightarrow \infty} E_{i}\right)=\lim _{i \rightarrow \infty} \mu\left(E_{i}\right)
$$

(v) (Continuity from the right) If $\left\{E_{i}\right\}$ is a decreasing sequence of sets, that is, if $E_{i} \supset E_{i+1}$ for each $i$, and if $\mu\left(E_{i_{0}}\right)<\infty$ for some $i_{0}$, then

$$
\mu\left(\bigcap_{i=1}^{\infty} E_{i}\right)=\mu\left(\lim _{i \rightarrow \infty} E_{i}\right)=\lim _{i \rightarrow \infty} \mu\left(E_{i}\right)
$$

(vi)

$$
\mu\left(\liminf _{i \rightarrow \infty} E_{i}\right) \leq \liminf _{i \rightarrow \infty} \mu\left(E_{i}\right)
$$

(vii) If

$$
\mu\left(\bigcup_{i=i_{0}}^{\infty} E_{i}\right)<\infty
$$

for some positive integer $i_{0}$, then

$$
\mu\left(\limsup _{i \rightarrow \infty} E_{i}\right) \geq \limsup _{i \rightarrow \infty} \mu\left(E_{i}\right)
$$

Proof. Only (i) and (iii) have not been established in Corollary 4.12. For (i), observe that if $E_{1} \subset E_{2}$, then $\mu\left(E_{2}\right)=\mu\left(E_{1}\right)+\mu\left(E_{2}-E_{1}\right) \geq \mu\left(E_{1}\right)$.
(iii) Refer to Lemma 4.7, to obtain a sequence of disjoint measurable sets $\left\{A_{i}\right\}$ such that $A_{i} \subset E_{i}$ and

$$
\bigcup_{i=1}^{\infty} E_{i}=\bigcup_{i=1}^{\infty} A_{i} .
$$

Then,

$$
\mu\left(\bigcup_{i=1}^{\infty} E_{i}\right)=\mu\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right) \leq \sum_{i=1}^{\infty} \mu\left(E_{i}\right)
$$

One property that is characteristic to an outer measure $\varphi$ but is not enjoyed by abstract measures in general is the following: if $\varphi(E)=0$, then $E$ is $\varphi$-measurable and consequently, so is every subset of $E$. Not all measures are complete, but this is not a crucial defect since every measure can easily be completed by enlarging its domain of definition to include all subsets of sets of measure zero.
4.50. Theorem. Suppose $(X, \mathcal{M}, \mu)$ is a measure space. Define $\overline{\mathcal{M}}=\{A \cup N$ : $A \in \mathcal{M}, N \subset B$ for some $B \in \mathcal{M}$ such that $\mu(B)=0\}$ and define $\bar{\mu}$ on $\overline{\mathcal{M}}$ by $\bar{\mu}(A \cup N)=\mu(A)$. Then, $\overline{\mathcal{M}}$ is a $\sigma$-algebra, $\bar{\mu}$ is a complete measure on $\overline{\mathcal{M}}$, and $(X, \overline{\mathcal{M}}, \bar{\mu})$ is a complete measure space. Moreover, $\bar{\mu}$ is the only complete measure on $\overline{\mathcal{M}}$ that is an extension of $\mu$.

Proof. It is easy to verify that $\overline{\mathcal{M}}$ is closed under countable unions since this is true for sets of measure zero. To show that $\overline{\mathcal{M}}$ is closed under complementation, note that with sets $A, N$, and $B$ as in the definition of $\overline{\mathcal{M}}$, it may be assumed that $A \cap N=\emptyset$ because $A \cup N=A \cup(N \backslash A)$ and $N \backslash A$ is a subset of a measurable set of measure zero, namely, $B \backslash A$. It can be readily verified that

$$
A \cup N=(A \cup B) \cap((\tilde{B} \cup N) \cup(A \cap B))
$$

and therefore

$$
\begin{aligned}
(A \cup N)^{\sim} & =(A \cup B)^{\sim} \cup((\tilde{B} \cup N) \cup(A \cap B))^{\sim} \\
& =(A \cup B)^{\sim} \cup\left((B \cap \tilde{N}) \cap(A \cap B)^{\sim}\right) \\
& =(A \cup B)^{\sim} \cup((B \backslash N) \backslash A \cap B) .
\end{aligned}
$$

Since $(A \cup B)^{\sim} \in \mathcal{M}$ and $(B \backslash N) \backslash A \cap B$ is a subset of a set of measure zero, it follows that $\overline{\mathcal{M}}$ is closed under complementation. Consequently, $\overline{\mathcal{M}}$ is a $\sigma$-algebra.

To show that the definition of $\bar{\mu}$ is unambiguous, suppose $A_{1} \cup N_{1}=A_{2} \cup N_{2}$ where $N_{i} \subset B_{i}, i=1,2$. Then $A_{1} \subset A_{2} \cup N_{2}$ and

$$
\bar{\mu}\left(A_{1} \cup N_{1}\right)=\mu\left(A_{1}\right) \leq \mu\left(A_{2}\right)+\mu\left(B_{2}\right)=\mu\left(A_{2}\right)=\bar{\mu}\left(A_{2} \cup N_{2}\right)
$$

Similarly, we have the opposite inequality. It is easily verified that $\bar{\mu}$ is complete since $\bar{\mu}(N)=\bar{\mu}(\emptyset \cup N)=\mu(\emptyset)=0$. Uniqueness is left as Exercise 2, Section 4.9.

## Exercises for Section 4.9

1. Let $\left\{\mu_{k}\right\}$ be a sequence of measures on a measure space such that $\mu_{k+1}(E) \geq$ $\mu_{k}(E)$ for each measurable set $E$. With $\mu$ defined as $\mu(E)=\lim _{k \rightarrow \infty} \mu_{k}(E)$, prove that $\mu$ is a measure.
2. Prove that the measure $\bar{\mu}$ introduced in Theorem 4.50 is a unique extension of $\mu$.
3. Let $\mu$ be finite Borel measure on $\mathbb{R}^{2}$. For fixed $r>0$, let $C_{x}=\{y:|y-x|=r\}$ and define $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ by $f(x)=\mu\left[C_{x}\right]$. Prove that $f$ is continuous at $x_{0}$ if and only if $\mu\left[C_{x_{0}}\right]=0$.
4. Let $\mu$ be finite Borel measure on $\mathbb{R}^{2}$. For fixed $r>0$, define $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ by $f(x)=\mu[B(x, r)]$. Prove that $f$ is continuous at $x_{0}$ if and only if $\mu\left[C_{x_{0}}\right]=0$.
5. This problem is set within the context of Theorem 4.50 of the text. With $\mu$ given as in Theorem 4.50, define an outer measure $\mu^{*}$ on all subsets of $X$ in the following way: For an arbitrary set $A \subset X$ let

$$
\mu^{*}(A):=\inf \left\{\sum_{i=1}^{\infty} \mu\left(E_{i}\right)\right\}
$$

where the infimum is taken over all countable collections $\left\{E_{i}\right\}$ such that

$$
A \subset \bigcup_{i=1}^{\infty} E_{i}, \quad E_{i} \in \mathcal{M}
$$

Prove that $\overline{\mathcal{M}}=\mathcal{M}^{*}$ where $\mathcal{M}^{*}$ denotes the $\sigma$-algebra of $\mu^{*}$-measurable sets and that $\bar{\mu}=\mu^{*}$ on $\overline{\mathcal{M}}$.
6. In an abstract measure space $(X, \mathcal{M}, \mu)$, if $\left\{A_{i}\right\}$ is a countable disjoint family of sets in $\mathcal{M}$, we know that

$$
\mu\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right)
$$

Prove that converse is essentially true. That is, under the assumption that $\mu(X)<\infty$, prove that if $\left\{A_{i}\right\}$ is a countable family of sets in $\mathcal{M}$ with the property that

$$
\mu\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right)
$$

then $\mu\left(A_{i} \cap A_{j}\right)=0$ whenever $i \neq j$.
7. Recall that an algebra in a space $X$ is a nonempty collection of subsets of $X$ that is closed under the operations of finite unions and complements. Also recall that a measure on an algebra, $\mathcal{A}$, is a function $\mu: \mathcal{A} \rightarrow[0, \infty]$ satisfying the properties
(i) $\mu(\emptyset)=0$,
(ii) if $\left\{A_{i}\right\}$ is a disjoint sequence of sets in $\mathcal{A}$ whose union is also in $\mathcal{A}$, then

$$
\mu\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right) .
$$

Finally, recall that a measure $\mu$ on an algebra $\mathcal{A}$ generates a set function $\mu^{*}$ defined on all subsets of $X$ in the following way: for each $E \subset X$, let

$$
\begin{equation*}
\mu^{*}(E):=\inf \left\{\sum_{i=1}^{\infty} \mu\left(A_{i}\right)\right\} \tag{4.38}
\end{equation*}
$$

where the infimum is taken over countable collections $\left\{A_{i}\right\}$ such that

$$
E \subset \bigcup_{i=1}^{\infty} A_{i}, \quad A_{i} \in \mathcal{A}
$$

Assuming that $\mu(X)<\infty$, prove that $\mu^{*}$ is a regular outer measure.
8. Give an example of two $\sigma$-algebras in a set $X$ whose union is not an algebra.
9. Prove that if the union of two $\sigma$-algebras is an algebra, then it is necessarily a $\sigma$-algebra.
10. Let $\varphi$ be an outer measure on a set $X$ and let $\mathcal{M}$ denote the $\sigma$-algebra of $\varphi$ measurable sets. Let $\mu$ denote the measure defined by $\mu(E)=\varphi(E)$ whenever $E \in \mathcal{M}$; that is, $\mu$ is the restriction of $\varphi$ to $\mathcal{M}$. Since, in particular, $\mathcal{M}$ is an algebra we know that $\mu$ generates an outer measure $\mu^{*}$. Prove:
(a) $\mu^{*}(E) \geq \varphi(E)$ whenever $E \in \mathcal{M}$
(b) $\mu^{*}(A)=\varphi(A)$ for $A \subset X$ if and only if there exists $E \in \mathcal{M}$ such that $E \supset A$ and $\varphi(E)=\varphi(A)$.
(c) $\mu^{*}(A)=\varphi(A)$ for all $A \subset X$ if $\varphi$ is regular.

### 4.10. Regular Outer Measures

In any context, the ability to approximate a complex entity by a simpler one is very important. The following result is one of many such approximations that occur in measure theory; it states that for outer measures with rather general properties, it is possible to approximate Borel sets by both open and closed sets. Note the strong parallel to similar results for Lebesgue measure and Hausdorff measure; see Theorems 4.25 and 4.37 along with Exercise 9, Section 4.3.
4.51. Theorem. If $\varphi$ is a regular outer measure on $X$, then
(i) If $A_{1} \subset A_{2} \subset \ldots$ is an increasing sequence of arbitrary sets, then

$$
\varphi\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\lim _{i \rightarrow \infty} \varphi\left(A_{i}\right)
$$

(ii) If $A \cup B$ is $\varphi$-measurable, $\varphi(A)<\infty, \varphi(B)<\infty$ and $\varphi(A \cup B)=\varphi(A)+\varphi(B)$, then both $A$ and $B$ are $\varphi$-measurable.

Proof. (i): Choose $\varphi$-measurable sets $C_{i} \supset A_{i}$ with $\varphi\left(C_{i}\right)=\varphi\left(A_{i}\right)$. The $\varphi$-measurable sets

$$
B_{i}:=\bigcap_{j=i}^{\infty} C_{j}
$$

form an ascending sequence that satisfy the conditions $A_{i} \subset B_{i} \subset C_{i}$ as well as

$$
\varphi\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \varphi\left(\bigcup_{i=1}^{\infty} B_{i}\right)=\lim _{i \rightarrow \infty} \varphi\left(B_{i}\right) \leq \lim _{i \rightarrow \infty} \varphi\left(C_{i}\right)=\lim _{i \rightarrow \infty} \varphi\left(A_{i}\right)
$$

Hence, it follows that

$$
\varphi\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \lim _{i \rightarrow \infty} \varphi\left(A_{i}\right)
$$

The opposite inequality is immediate since

$$
\varphi\left(\bigcup_{i=1}^{\infty} A_{i}\right) \geq \varphi\left(A_{k}\right) \text { for each } k \in \mathbb{N}
$$

(ii): Choose a $\varphi$-measurable set $C^{\prime} \supset A$ such that $\varphi\left(C^{\prime}\right)=\varphi(A)$. Then, with $C:=C^{\prime} \cap(A \cup B)$, we have a $\varphi$-measurable set $C$ with $A \subset C \subset A \cup B$ and
$\varphi(C)=\varphi(A)$. Note that

$$
\begin{equation*}
\varphi(B \cap C)=0 \tag{4.39}
\end{equation*}
$$

because the $\varphi$-measurability of $C$ implies

$$
\varphi(B)=\varphi(B \cap C)+\varphi(B \backslash C)
$$

and

$$
\begin{aligned}
\varphi(C)+\varphi(B) & =\varphi(A)+\varphi(B) \\
& =\varphi(A \cup B) \\
& =\varphi((A \cup B) \cap C)+\varphi((A \cup B) \backslash C) \\
& =\varphi(C)+\varphi(B \backslash C) \\
& =\varphi(C)+\varphi(B)-\varphi(B \cap C)
\end{aligned}
$$

This implies that $\varphi(B \cap C)=0$ because $\varphi(B)+\varphi(C)<\infty$. Since $C \subset A \cup B$ we have
$C \backslash A \subset B$ which leads to $(C \backslash A) \subset B \cap C$. Then (4.39) implies $\varphi(C \backslash A)=0$ which yields the $\varphi$-measurability of $A$ since $A=C \backslash(C \backslash A)$. $B$ is also $\varphi$-measurable, since the roles of $A$ and $B$ are interchangeable.
4.52. Theorem. Suppose $\varphi$ is an outer measure on a metric space $X$ whose measurable sets contain the Borel sets; that is, $\varphi$ is a Borel outer measure. Then, for each Borel set $B \subset X$ with $\varphi(B)<\infty$ and each $\varepsilon>0$, there exists $a$ closed set $F \subset B$ such that

$$
\varphi(B \backslash F)<\varepsilon .
$$

Furthermore, suppose

$$
B \subset \bigcup_{i=1}^{\infty} V_{i}
$$

where each $V_{i}$ is an open set with $\varphi\left(V_{i}\right)<\infty$. Then, for each $\varepsilon>0$, there is an open set $W \supset B$ such that

$$
\varphi(W \backslash B)<\varepsilon
$$

Proof. For the proof of the first part, select a Borel set $B$ with $\varphi(B)<\infty$ and define a set function $\mu$ by

$$
\begin{equation*}
\mu(A)=\varphi(A \cap B) \tag{4.40}
\end{equation*}
$$

whenever $A \subset X$. It is easy to verify that $\mu$ is an outer measure on $X$ whose measurable sets include all $\varphi$-measurable sets (see Exercise 3, Section 4.2) and thus all open sets. The outer measure $\mu$ is introduced merely to allow us to work with an outer measure for which $\mu(X)<\infty$.

Let $\mathcal{D}$ be the family of all $\mu$-measurable sets $A \subset X$ with the following property: For each $\varepsilon>0$, there is a closed set $F \subset A$ such that $\mu(A \backslash F)<\varepsilon$. The first part of the Theorem will be established by proving that $\mathcal{D}$ contains all Borel sets. Obviously, $\mathcal{D}$ contains all closed sets. It also contains all open sets. Indeed, if $U$ is an open set, then the closed sets

$$
F_{i}=\{x: \boldsymbol{d}(x, \tilde{U}) \geq 1 / i\}
$$

have the property that $F_{1} \subset F_{2} \subset \ldots$ and

$$
U=\bigcup_{i=1}^{\infty} F_{i}
$$

and therefore that

$$
\bigcap_{i=1}^{\infty}\left(U \backslash F_{i}\right)=\emptyset .
$$

Therefore, since $\mu(X)<\infty$, Corollary 4.12 (iv) yields

$$
\lim _{i \rightarrow \infty} \mu\left(U \backslash F_{i}\right)=0
$$

which shows that $\mathcal{D}$ contains all open sets $U$.
Since $\mathcal{D}$ contains all open and closed sets, according to Proposition 4.17, we need only show that $\mathcal{D}$ is closed under countable unions and countable intersections to conclude that it also contains all Borel sets. For this purpose, suppose $\left\{A_{i}\right\}$ is a sequence of sets in $\mathcal{D}$ and for given $\varepsilon>0$, choose closed sets $C_{i} \subset A_{i}$ with $\mu\left(A_{i} \backslash C_{i}\right)<\varepsilon / 2^{i}$. Since

$$
\bigcap_{i=1}^{\infty} A_{i} \backslash \bigcap_{i=1}^{\infty} C_{i} \subset \bigcup_{i=1}^{\infty}\left(A_{i} \backslash C_{i}\right)
$$

and

$$
\bigcup_{i=1}^{\infty} A_{i} \backslash \bigcup_{i=1}^{\infty} C_{i} \subset \bigcup_{i=1}^{\infty}\left(A_{i} \backslash C_{i}\right)
$$

it follows that

$$
\begin{equation*}
\mu\left[\bigcap_{i=1}^{\infty} A_{i} \backslash \bigcap_{i=1}^{\infty} C_{i}\right] \leq \mu\left[\bigcup_{i=1}^{\infty}\left(A_{i} \backslash C_{i}\right)\right]<\sum_{i=1}^{\infty} \frac{\varepsilon}{2^{i}}=\varepsilon \tag{4.41}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \mu\left[\bigcup_{i=1}^{\infty} A_{i} \backslash \bigcup_{i=1}^{k} C_{i}\right]=\mu\left[\bigcup_{i=1}^{\infty} A_{i} \backslash \bigcup_{i=1}^{\infty} C_{i}\right] \leq \mu\left[\bigcup_{i=1}^{\infty}\left(A_{i} \backslash C_{i}\right)\right]<\varepsilon \tag{4.42}
\end{equation*}
$$

Consequently, there exists a positive integer $k$ such that

$$
\begin{equation*}
\mu\left[\bigcup_{i=1}^{\infty} A_{i} \backslash \bigcup_{i=1}^{k} C_{i}\right]<\varepsilon . \tag{4.43}
\end{equation*}
$$

We have used the fact that $\cup_{i=1}^{\infty} A_{i}$ and $\cap_{i=1}^{\infty} A_{i}$ are $\mu$-measurable, and in (4.42), we again have used (iv) of Corollary 4.12. Since the sets $\cap_{i=1}^{\infty} C_{i}$ and $\cup_{i=1}^{k} C_{i}$ are closed
subsets of $\cap{ }_{i=1}^{\infty} A_{i}$ and $\cup_{i=1}^{\infty} A_{i}$ respectively, it follows from (4.41) and (4.43) that $\mathcal{D}$ is closed under the operations of countable unions and intersections.

To prove the second part of the Theorem, consider the Borel sets $V_{i} \backslash B$ and use the first part to find closed sets $C_{i} \subset\left(V_{i} \backslash B\right)$ such that

$$
\varphi\left[\left(V_{i} \backslash C_{i}\right) \backslash B\right]=\varphi\left[\left(V_{i} \backslash B\right) \backslash C_{i}\right]<\frac{\varepsilon}{2^{i}}
$$

For the desired set $W$ in the statement of the Theorem, let $W=\cup_{i=1}^{\infty}\left(V_{i} \backslash C_{i}\right)$ and observe that

$$
\varphi(W \backslash B) \leq \sum_{i=1}^{\infty} \varphi\left[\left(V_{i} \backslash C_{i}\right) \backslash B\right]<\sum_{i=1}^{\infty} \frac{\varepsilon}{2^{i}}=\varepsilon
$$

Moreover, since $B \cap V_{i} \subset V_{i} \backslash C_{i}$, we have

$$
B=\bigcup_{i=1}^{\infty}\left(B \cap V_{i}\right) \subset \bigcup_{i=1}^{\infty}\left(V_{i} \backslash C_{i}\right)=W
$$

4.53. Corollary. If two finite Borel outer measures agree on all open (or closed) sets, then they agree on all Borel sets. In particular, in $\mathbb{R}$, if they agree on all half-open intervals, then they agree on all Borel sets.
4.54. Remark. The preceding theorem applies directly to any Carathéodory outer measure since its measurable sets contain the Borel sets. In particular, the result applies to both Lebesgue- Stieltjes measure $\lambda_{f}^{*}$ and Lebesgue measure and thereby furnishes an alternate proof to Theorem 4.25.
4.55. Remark. In order to underscore the importance of Theorem 4.52, let us return to Theorem 4.33. There we are given a Borel outer measure $\mu$ with $\mu(\mathbb{R})<\infty$. Then define a function $f$ by

$$
f(x):=\mu((-\infty, x])
$$

and observe that $f$ is nondecreasing and right-continuous. Consequently, $f$ produces a Lebesgue- Stieltjes measure $\lambda_{f}^{*}$ with the property that

$$
\lambda_{f}^{*}((a, b])=f(b)-f(a)
$$

for each half-open interval. However, it is clear from the definition of $f$ that $\mu$ also enjoys the same property:

$$
\mu((a, b])=f(b)-f(a)
$$

Thus, $\mu$ and $\lambda_{f}^{*}$ agree on all half-open intervals and therefore they agree on all open sets, since any open set is the disjoint union of half-open intervals. Hence, from Corollary 4.53, they agree on all Borel sets. This allows us to conclude that there is unique correspondance between nondecreasing, right-continuous functions and finite Borel measures on $\mathbb{R}$.

It is natural to ask whether the previous theorem remains true if $B$ is only assumed to be $\varphi$-measurable rather than being a Borel set. In general the answer is no, but it is true if $\varphi$ is assumed to be a Borel regular outer measure. To see this, observe that if $\varphi$ is a Borel regular outer measure and $A$ is a $\varphi$-measurable set with $\varphi(A)<\infty$, then there exist Borel sets $B_{1}$ and $B_{2}$ such that

$$
\begin{equation*}
B_{2} \subset A \subset B_{1} \quad \text { and } \quad \varphi\left(B_{1} \backslash B_{2}\right)=0 \tag{4.44}
\end{equation*}
$$

Proof. For this, first choose a Borel set $B_{1} \supset A$ with $\varphi\left(B_{1}\right)=\varphi(A)$. Then choose a Borel set $D \supset B_{1} \backslash A$ such that $\varphi(D)=\varphi\left(B_{1} \backslash A\right)$. Note that since $A$ and $B_{1}$ are $\varphi$-measurable, we have $\varphi\left(B_{1} \backslash A\right)=\varphi\left(B_{1}\right)-\varphi(A)=0$. Now take $B_{2}=B_{1} \backslash D$. Thus, we have the following corollary.
4.56. Corollary. In the previous theorem, if $\varphi$ is assumed to be a Borel regular outer measure, then the conclusions remain valid if the phrase"for each Borel set B" is replaced by "for each $\varphi$-measurable set B."

Although not all Carathéodory outer measures are Borel regular, the following theorems show that they do agree with Borel regular outer measures on the Borel sets.
4.57. Theorem. Let $\varphi$ be a Carathéodory outer measure. For each set $A \subset X$, define

$$
\begin{equation*}
\psi(A)=\inf \{\varphi(B): B \supset A, B \text { a Borel set }\} \tag{4.45}
\end{equation*}
$$

Then $\psi$ is a Borel regular outer measure on $X$, which agrees with $\varphi$ on all Borel sets.

Proof. We leave it as an easy exercise (Exercise 4, Section 4.2) to show that $\psi$ is an outer measure on $X$. To show that all Borel sets are $\psi$-measurable, suppose $D \subset X$ is a Borel set. Then, by Definition 4.3, we must show that

$$
\begin{equation*}
\psi(A) \geq \psi(A \cap D)+\psi(A \backslash D) \tag{4.46}
\end{equation*}
$$

whenever $A \subset X$. For this we may as well assume $\psi(A)<\infty$. For $\varepsilon>0$, choose a Borel set $B \supset A$ such that $\varphi(B)<\psi(A)+\varepsilon$. Then, since $\varphi$ is a Borel outer measure (Theorem 4.18), we have

$$
\begin{aligned}
\varepsilon+\psi(A) \geq \varphi(B) & \geq \varphi(B \cap D)+\varphi(B \backslash D) \\
& \geq \psi(A \cap D)+\psi(A \backslash D)
\end{aligned}
$$

which establishes (4.46) since $\varepsilon$ is arbitrary. Also, if $B$ is a Borel set, we claim that $\psi(B)=\varphi(B)$. Half the claim is obvious because $\psi(B) \leq \varphi(B)$ by definition. As for the opposite inequality, choose a sequence of Borel sets $D_{i} \subset X$ with $D_{i} \supset B$ and $\lim _{i \rightarrow \infty} \varphi\left(D_{i}\right)=\psi(B)$. Then, with $D=\liminf _{i \rightarrow \infty} D_{i}$, we have by Corollary 4.12 (v)

$$
\varphi(B) \leq \varphi(D) \leq \liminf _{i \rightarrow \infty} \varphi\left(D_{i}\right)=\psi(B)
$$

which establishes the claim. Finally, since $\varphi$ and $\psi$ agree on Borel sets, we have for arbitrary $A \subset X$,

$$
\begin{aligned}
\psi(A) & =\inf \{\varphi(B): B \supset A, B \text { a Borel set }\} \\
& =\inf \{\psi(B): B \supset A, B \text { a Borel set }\}
\end{aligned}
$$

For each positive integer $i$, let $B_{i} \supset A$ be a Borel set with $\psi\left(B_{i}\right)<\psi(A)+1 / i$. Then

$$
B=\bigcap_{i=1}^{\infty} B_{i} \supset A
$$

is a Borel set with $\psi(B)=\psi(A)$, which shows that $\psi$ is Borel regular.

### 4.11. Outer Measures Generated by Measures

Thus far we have seen that with every outer measure there is an associated measure. This measure is defined by restricting the outer measure to its measurable sets. In this section, we consider the situation in reverse. It is shown that a measure defined on an abstract space generates an outer measure and that if this measure is $\sigma$-finite, the extension is unique. An important consequence of this development is that any finite Borel measure is necessarily regular.

We begin by describing a process by which a measure generates an outer measure. This method is reminiscent of the one used to define Lebesgue- Stieltjes measure. Actually, this method does not require the measure to be defined on a $\sigma$-algebra, but rather, only on an algebra of sets. We make this precise in the following Definition.
4.58. Definitions. An algebra in a space $X$ is defined as a nonempty collection of subsets of $X$ that is closed under the operations of finite unions and complements. Thus, the only difference between an algebra and a $\sigma$-algebra is that the latter is closed under countable unions. By a measure on an algebra, $\mathcal{A}$, we mean a function $\mu: \mathcal{A} \rightarrow[0, \infty]$ satisfying the properties
(i) $\mu(\emptyset)=0$,
(ii) if $\left\{A_{i}\right\}$ is a disjoint sequence of sets in $\mathcal{A}$ whose union is also in $\mathcal{A}$, then

$$
\mu\left(\bigcup_{i=1}^{\infty} A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right)
$$

Consequently, a measure on an algebra $\mathcal{A}$ is a measure (in the sense of Definition 4.47) if and only if $\mathcal{A}$ is a $\sigma$-algebra. A measure on $\mathcal{A}$ is called $\sigma$-finite if $X$ can be written

$$
X=\bigcup_{i=1}^{\infty} A_{i}
$$

with $A_{i} \in \mathcal{A}$ and $\mu\left(A_{i}\right)<\infty$.
A measure $\mu$ on an algebra $\mathcal{A}$ generates a set function $\mu^{*}$ defined on all subsets of $X$ in the following way: for each $E \subset X$, let

$$
\begin{equation*}
\mu^{*}(E):=\inf \left\{\sum_{i=1}^{\infty} \mu\left(A_{i}\right)\right\} \tag{4.47}
\end{equation*}
$$

where the infimum is taken over countable collections $\left\{A_{i}\right\}$ such that

$$
E \subset \bigcup_{i=1}^{\infty} A_{i}, \quad A_{i} \in \mathcal{A}
$$

Note that this definition is in the same spirit as that used to define Lebesgue measure or more generally, Lebesgue- Stieltjes measure.
4.59. Theorem. Let $\mu$ be a measure on an algebra $\mathcal{A}$ and let $\mu^{*}$ be the corresponding set function generated by $\mu$. Then
(i) $\mu^{*}$ is an outer measure,
(ii) $\mu^{*}$ is an extension of $\mu$; that is, $\mu^{*}(A)=\mu(A)$ whenever $A \in \mathcal{A}$.
(iii) each $A \in \mathcal{A}$ is $\mu^{*}$-measurable.

Proof. The proof of (i) is similar to showing that $\lambda^{*}$ is an outer measure (see the proof of Theorem 4.23) and is left as an exercise.
(ii) From the definition, $\mu^{*}(A) \leq \mu(A)$ whenever $A \in \mathcal{A}$. For the opposite inequality, consider $A \in \mathcal{A}$ and let $\left\{A_{i}\right\}$ be any sequence of sets in $\mathcal{A}$ with

$$
A \subset \bigcup_{i=1}^{\infty} A_{i} .
$$

Set

$$
B_{i}=A \cap A_{i} \backslash\left(A_{i-1} \cup A_{i-2} \cup \cdots \cup A_{1}\right) .
$$

These sets are disjoint. Furthermore, $B_{i} \in \mathcal{A}, B_{i} \subset A_{i}$, and $A=\cup_{i=1}^{\infty} B_{i}$. Hence, by the countable additivity of $\mu$,

$$
\mu(A)=\sum_{i=1}^{\infty} \mu\left(B_{i}\right) \leq \sum_{i=1}^{\infty} \mu\left(A_{i}\right)
$$

Since, by definition, the infimum of the right-side of this expression tends to $\mu^{*}(A)$, this shows that $\mu(A) \leq \mu^{*}(A)$.
(iii) For $A \in \mathcal{A}$, we must show that

$$
\mu^{*}(E) \geq \mu^{*}(E \cap A)+\mu^{*}(E \backslash A)
$$

whenever $E \subset X$. For this we may assume that $\mu^{*}(E)<\infty$. Given $\varepsilon>0$, there is a sequence of sets $\left\{A_{i}\right\}$ in $\mathcal{A}$ such that

$$
E \subset \bigcup_{i=1}^{\infty} A_{i} \quad \text { and } \quad \sum_{i=1}^{\infty} \mu\left(A_{i}\right)<\mu^{*}(E)+\varepsilon
$$

Since $\mu$ is additive on $\mathcal{A}$, we have

$$
\mu\left(A_{i}\right)=\mu\left(A_{i} \cap A\right)+\mu\left(A_{i} \backslash A\right)
$$

In view of the inclusions

$$
E \cap A \subset \bigcup_{i=1}^{\infty}\left(A_{i} \cap A\right) \quad \text { and } \quad E \backslash A \subset \bigcup_{i=1}^{\infty}\left(A_{i} \backslash A\right)
$$

we have

$$
\begin{aligned}
\mu^{*}(E)+\varepsilon & >\sum_{i=1}^{\infty} \mu\left(A_{i} \cap A\right)+\sum_{i=1}^{\infty} \mu\left(A_{i} \cap \tilde{A}\right) \\
& >\mu^{*}(E \cap A)+\mu^{*}(E \backslash A) .
\end{aligned}
$$

Since $\varepsilon$ is arbitrary, the desired result follows.
4.60. Example. Let us see how the previous result can be used to produce Lebesgue- Stieltjes measure. Let $\mathcal{A}$ be the algebra formed by including $\emptyset, \mathbb{R}$, all intervals of the form $(-\infty, a],(b,+\infty)$ along with all possible finite disjoint unions of these and intervals of the form $(a, b]$. Suppose that $f$ is a nondecreasing, rightcontinuous function and define $\mu$ on intervals $(a, b]$ in $\mathcal{A}$ by

$$
\mu((a, b])=f(b)-f(a)
$$

and then extend $\mu$ to all elements of $\mathcal{A}$ by additivity. Then we see that the outer measure $\mu^{*}$ generated by $\mu$ using (4.47) agrees with the definition of LebesgueStieltjes measure defined by (4.23). Our previous result states that $\mu^{*}(A)=\mu(A)$ for all $A \in \mathcal{A}$, which agrees with Theorem 4.32.
4.61. Remark. In the previous example, the right-continuity of $f$ is needed to ensure that $\mu$ is in fact a measure on $\mathcal{A}$. For example, if

$$
f(x):= \begin{cases}0 & x \leq 0 \\ 1 & x>0\end{cases}
$$

then $\mu((0,1])=1$. But $(0,1]=\bigcup_{k=1}^{\infty}\left(\frac{1}{k+1}, \frac{1}{k}\right]$ and

$$
\begin{aligned}
\mu\left(\bigcup_{k=1}^{\infty}\left(\frac{1}{k+1}, \frac{1}{k}\right]\right) & =\sum_{k=1}^{\infty} \mu\left(\left(\frac{1}{k+1}, \frac{1}{k}\right]\right) \\
& =0
\end{aligned}
$$

which shows that $\mu$ is not a measure.

Next is the main result of this section, which in addition to restating the results of Theorem 4.59, ensures that the outer measure generated by $\mu$ is unique.
4.62. Theorem (Carathéodory-Hahn Extension Theorem). Let $\mu$ be a measure on an algebra $\mathcal{A}$, let $\mu^{*}$ be the outer measure generated by $\mu$, and let $\mathcal{A}^{*}$ be the $\sigma$ algebra of $\mu^{*}$-measurable sets.
(i) Then $\mathcal{A}^{*} \supset \mathcal{A}$ and $\mu^{*}=\mu$ on $\mathcal{A}$,
(ii) Let $\mathcal{M}$ be a $\sigma$-algebra with $\mathcal{A} \subset \mathcal{M} \subset \mathcal{A}^{*}$ and suppose $\nu$ is a measure on $\mathcal{M}$ that agrees with $\mu$ on $\mathcal{A}$. Then $\nu=\mu^{*}$ on $\mathcal{M}$ provided that $\mu$ is $\sigma$-finite.

Proof. As noted above, (i) is a restatement of Theorem 4.59.
(ii) Given any $E \in \mathcal{M}$ note that $\nu(E) \leq \mu^{*}(E)$ since if $\left\{A_{i}\right\}$ is any countable collection in $\mathcal{A}$ whose union contains $E$, then

$$
\nu(E) \leq \nu\left(\bigcup_{i=1}^{\infty} A_{i}\right) \leq \sum_{i=1}^{\infty} \nu\left(A_{i}\right)=\sum_{i=1}^{\infty} \mu\left(A_{i}\right)
$$

To prove equality let $A \in \mathcal{A}$ with $\mu(A)<\infty$. Then we have

$$
\begin{equation*}
\nu(E)+\nu(A \backslash E)=\nu(A)=\mu^{*}(A)=\mu^{*}(E)+\mu^{*}(A \backslash E) \tag{4.48}
\end{equation*}
$$

Note that $A \backslash E \in \mathcal{M}$, and therefore $\nu(A \backslash E) \leq \mu^{*}(A \backslash E)$ from what we have just proved. Since all terms in (4.48) are finite we deduce that

$$
\nu(A \cap E)=\mu^{*}(A \cap E)
$$

whenever $A \in \mathcal{A}$ with $\mu(A)<\infty$. Since $\mu$ is $\sigma$-finite, there exist $A_{i} \in \mathcal{A}$ such that

$$
X=\bigcup_{i=1}^{\infty} A_{i}
$$

with $\mu\left(A_{i}\right)<\infty$ for each $i$. We may assume that the $A_{i}$ are disjoint (Lemma 4.7) and therefore

$$
\nu(E)=\sum_{i=1}^{\infty} \nu\left(E \cap A_{i}\right)=\sum_{i=1}^{\infty} \mu^{*}\left(E \cap A_{i}\right)=\mu^{*}(E)
$$

Let us consider a special case of this result, namely, the situation in which $\mathcal{A}$ is the family of Borel sets in a metric space $X$. If $\mu$ is a finite measure defined on the Borel sets, the previous result states that the outer measure, $\mu^{*}$, generated by $\mu$ agrees with $\mu$ on the Borel sets. Theorem 4.52 asserts that $\mu^{*}$ enjoys certain regularity properties. Since $\mu$ and $\mu^{*}$ agree on Borel sets, it follows that $\mu$ also enjoys these regularity properties. This implies the remarkable fact that any finite Borel measure is automatically regular. We state this as our next result.
4.63. Theorem. Suppose $(X, \mathcal{M}, \mu)$ is a measure space where $X$ is a metric space and $\mu$ is a finite Borel measure (that is, $\mathcal{M}$ denotes the Borel sets of $X$ and $\mu(X)<\infty)$. Then for each $\varepsilon>0$ and each Borel set $B$, there is an open set $U$ and a closed set $F$ such that $F \subset B \subset U, \mu(B \backslash F)<\varepsilon$ and $\mu(U \backslash B)<\varepsilon$.

In case $\mu$ is a measure defined on a $\sigma$-algebra $\mathcal{M}$ rather than on an algebra $\mathcal{A}$, there is another method for generating an outer measure. In this situation, we define $\mu^{* *}$ on an arbitrary set $E \subset X$ by

$$
\begin{equation*}
\mu^{* *}(E)=\inf \{\mu(B): B \supset E, B \in \mathcal{M}\} \tag{4.49}
\end{equation*}
$$

We have the following result.
4.64. Theorem. Consider a measure space $(X, \mathcal{M}, \mu)$. The set function $\mu^{* *}$ defined above is an outer measure on $X$. Moreover, $\mu^{* *}$ is a regular outer measure and $\mu(B)=\mu^{* *}(B)$ for each $B \in \mathcal{M}$.

Proof. The proof proceeds exactly as in Theorem 4.57. One need only replace each reference to a Borel set in that proof with $\mathcal{M}$-measurable set.
4.65. Theorem. Suppose $(X, \mathcal{M}, \mu)$ is a measure space and let $\mu^{*}$ and $\mu^{* *}$ be the outer measures generated by $\mu$ as described in (4.47) and (4.49), respectively. Then, for each $E \subset X$ with $\mu(E)<\infty$ there exists $B \in \mathcal{M}$ such that $B \supset E$,

$$
\mu(B)=\mu^{*}(B)=\mu^{*}(E)=\mu^{* *}(E)
$$

Proof. We will show that for any $E \subset X$ there exists $B \in \mathcal{M}$ such that $B \supset E$ and $\mu(B)=\mu^{*}(B)=\mu^{*}(E)$. From the previous result, it will then follow that $\mu^{*}(E)=\mu^{* *}(E)$.

Note that for each $\varepsilon>0$, and any set $E$, there exists a sequence $\left\{A_{i}\right\} \in \mathcal{M}$ such that

$$
E \subset \bigcup_{i=1}^{\infty} A_{i} \quad \text { and } \quad \sum_{i=1}^{\infty} \mu\left(A_{i}\right) \leq \mu^{*}(E)+\varepsilon
$$

Setting $A=\cup A_{i}$, we have

$$
\mu(A)<\mu^{*}(E)+\varepsilon
$$

For each positive integer $k$, use this observation with $\varepsilon=1 / k$ to obtain a set $A_{k} \in \mathcal{M}$ such that $A_{k} \supset E$ and $\mu\left(A_{k}\right)<\mu^{*}(E)+1 / k$. Let

$$
B=\bigcap_{k=1}^{\infty} A_{k}
$$

Then $B \in \mathcal{M}$ and since $E \subset B \subset A_{k}$, we have

$$
\mu^{*}(E) \leq \mu^{*}(B) \leq \mu(B) \leq \mu\left(A_{k}\right)<\mu^{*}(E)+1 / k
$$

Since $k$ is arbitrary, it follows that $\mu(B)=\mu^{*}(B)=\mu^{*}(E)$.

## CHAPTER 5

## Measurable Functions

### 5.1. Elementary Properties of Measurable Functions

The class of measurable functions will play a critical role in the theory of integration. It is shown that this class remains closed under the usual elementary operations, although special care must be taken in the case of composition of functions. The main results of this chapter are the theorems of Egoroff and Lusin. Roughly, they state that pointwise convergence of a sequence of measurable functions is "nearly" uniform convergence and that a measurable function is "nearly" continuous.

Throughout this chapter, we will consider an abstract measure space ( $X, \mathcal{M}, \mu$ ), where $\mu$ is a measure defined on the $\sigma$-algebra $\mathcal{M}$. Virtually all the material in this first section depends only on the $\sigma$-algebra and not on the measure $\mu$. This is a reflection of the fact that the elementary properties of measurable functions are settheoretic and are not related to $\mu$. Also, we will consider functions $f: X \rightarrow \overline{\mathbb{R}}$, where $\overline{\mathbb{R}}=\mathbb{R} \cup\{-\infty\} \cup\{+\infty\}$ is the set of extended real numbers. For convenience, we will write $\infty$ for $+\infty$. Arithmetic operations on $\overline{\mathbb{R}}$ are subject to the following conventions. For $x \in \mathbb{R}$, we define

$$
x+( \pm \infty)=( \pm \infty)+x= \pm \infty
$$

and

$$
( \pm \infty)+( \pm \infty)= \pm \infty,( \pm \infty)-(\mp \infty)= \pm \infty
$$

but

$$
( \pm \infty)+(\mp \infty), \text { and }( \pm \infty)-( \pm \infty)
$$

are undefined. Also, for the operation of multiplication, we define

$$
x( \pm \infty)=( \pm \infty) x= \begin{cases} \pm \infty, & x>0 \\ 0, & x=0 \\ \mp \infty, & x<0\end{cases}
$$

for each $x \in \mathbb{R}$ and let

$$
( \pm \infty) \cdot( \pm \infty)=+\infty \quad \text { and } \quad( \pm \infty) \cdot(\mp \infty)=-\infty
$$

The operations

$$
\frac{\infty}{-\infty}, \frac{-\infty}{\infty}, \frac{\infty}{\infty} \text { and } \frac{-\infty}{-\infty}
$$

are undefined.
We endow $\overline{\mathbb{R}}$ with a topology called the order topology in the following manner. For each $a \in \mathbb{R}$ let

$$
L_{a}=\overline{\mathbb{R}} \cap\{x: x<a\}=[-\infty, a) \quad \text { and } \quad R_{a}=\overline{\mathbb{R}} \cap\{x: x>a\}=(a, \infty] .
$$

The collection $\mathcal{S}=\left\{L_{a}: a \in \mathbb{R}\right\} \cup\left\{R_{a}: a \in \mathbb{R}\right\}$ is taken as a subbase for this topology. A base for the topology is given by

$$
\mathcal{S} \cup\left\{R_{a} \cap L_{b}: a, b \in \mathbb{R}, a<b\right\} .
$$

Observe that the topology on $\mathbb{R}$ induced by the order topology on $\overline{\mathbb{R}}$ is precisely the usual topology on $\mathbb{R}$.

Suppose $X$ and $Y$ are topological spaces. Recall that a mapping $f: X \rightarrow Y$ is continuous if and only if $f^{-1}(U)$ is open whenever $U \subset Y$ is open. We define a measurable mapping analogously.
5.1. Definitions. Suppose $(X, \mathcal{M})$ and $(Y, \mathcal{N})$ are measure spaces. A mapping $f: X \rightarrow Y$ is called measurable with respect to $\mathcal{M}$ and $\mathcal{N}$ if

$$
\begin{equation*}
f^{-1}(E) \in \mathcal{M} \text { whenever } E \in \mathcal{N} \tag{5.1}
\end{equation*}
$$

If there is no danger of confusion, reference to $\mathcal{M}$ and $\mathcal{N}$ will be omitted, and we will simply use the term "measurable mapping."

In case $Y$ is a topological space, a restriction is placed on $\mathcal{N}$. In this case it is always assumed that $\mathcal{N}$ is the $\sigma$-algebra of Borel sets $\mathcal{B}$. Thus, in this situation, a mapping $(X, \mathcal{M}) \xrightarrow{f}(Y, \mathcal{B})$ is measurable if

$$
\begin{equation*}
f^{-1}(E) \in \mathcal{M} \text { whenever } E \in \mathcal{B} \tag{5.2}
\end{equation*}
$$

The reason for imposing this condition is to ensure that continuous mappings will be measurable. That is, if both $X$ and $Y$ are topological spaces, $X \xrightarrow{f} Y$ is continuous and $\mathcal{M}$ contains the Borel sets of $X$, then $f$ is measurable, since $f^{-1}(E) \in \mathcal{M}$ whenever $E$ is a Borel set, see Exercise 1, Section 5.1. One of the most important situations is when $Y$ is taken as $\overline{\mathbb{R}}$ (endowed with the order topology) and $(X, \mathcal{M})$ is a topological space with $\mathcal{M}$ the collection of Borel sets. Then $f$ is called a Borel measurable function. Another important example of this is when $X=\mathbb{R}^{n}, \mathcal{M}$ is the class of Lebesgue measurable sets and $Y=\overline{\mathbb{R}}$. Here, it is required that $f^{-1}(E)$ is Lebesgue measurable whenever $E \subset \overline{\mathbb{R}}$ is Borel, in which case $f$ is called a Lebesgue measurable function. The definitions imply that $E$ is a measurable set if and only if $\chi_{E}$ is a measurable function.

If the mapping $(X, \mathcal{M}) \xrightarrow{f}(Y, \mathcal{B})$ is measurable, where $\mathcal{B}$ is the $\sigma$-algebra of Borel sets, then we can make the following observation which will be useful in the development. Define

$$
\begin{equation*}
\Sigma=\left\{E: E \subset Y \text { and } f^{-1}(E) \in \mathcal{M}\right\} \tag{5.3}
\end{equation*}
$$

Note that $\Sigma$ is closed under countable unions. It is also closed under complementation since

$$
\begin{equation*}
f^{-1}\left(E^{\sim}\right)=\left[f^{-1}(E)\right]^{\sim} \in \mathcal{M} \tag{5.4}
\end{equation*}
$$

for $E \in \Sigma$ and thus, $\Sigma$ is a $\sigma$-algebra.
In view of (5.3) and (5.4), note that a continuous mapping is a Borel measurable function (Exercise 1, Section 5.1 ).

In case $f: X \rightarrow \overline{\mathbb{R}}$, it will be convenient to characterize measurability in terms of the sets $X \cap\{x: f(x)>a\}$ for $a \in \mathbb{R}$. To simplify notation, we simply write $\{f>a\}$ to denote these sets. The sets $\{f>a\}$ are called the superlevel sets of $f$. The behavior of a function $f$ is to a large extent reflected in the properties of its superlevel sets. For example, if $f$ is a continuous function on a metric space $X$, then $\{f>a\}$ is an open set for each real number $a$. If the function is nicer, then we should expect better behavior of the superlevel sets. Indeed, if $f$ is an infinitely differentiable function $f$ defined on $\mathbb{R}^{n}$ with nonvanishing gradient, then not only is each $\{f>a\}$ an open set, but an application of the Implicit Function Theorem shows that its boundary is a smooth manifold of dimension $n-1$ as well.

We begin by showing that the definition of an $\overline{\mathbb{R}}$-valued measurable function could just as well be stated in terms of its level sets.
5.2. Theorem. Let $f: X \rightarrow \overline{\mathbb{R}}$ where $(X, \mathcal{M})$ is a measure space. The following conditions are equivalent:
(i) $f$ is measurable.
(ii) $\{f>a\} \in \mathcal{M}$ for each $a \in \mathbb{R}$.
(iii) $\{f \geq a\} \in \mathcal{M}$ for each $a \in \mathbb{R}$.
(iv) $\{f<a\} \in \mathcal{M}$ for each $a \in \mathbb{R}$.
(v) $\{f \leq a\} \in \mathcal{M}$ for each $a \in \mathbb{R}$.

Proof. (i) implies (ii) by definition since $\{f>a\}=f^{-1}((a, \infty])$ and $(a, \infty]$ is open in the order topology. In view of $\{f \geq a\}=\cap_{k=1}^{\infty}\{f>a-1 / k\}$, (ii) implies (iii). The set $\{f<a\}$ is the complement of $\{f \geq a\}$, thus establishing the next implication. Similar to the proof of the first implication, we have $\{f \leq a\}=$ $\cap_{k=1}^{\infty}\{f<a+1 / k\}$, which shows that (iv) implies (v). For the proof that (v) implies
(i), in view of (5.3) and (5.4) with $Y=\overline{\mathbb{R}}$, it is sufficient to show that $f^{-1}(U) \in \mathcal{M}$ whenever $U \subset \overline{\mathbb{R}}$ is open. Since $f^{-1}$ preserves unions and intersections and $U$ can be written as a countable union of elements of the base, we need only consider $f^{-1}(J)$ where $J$ assumes the form $J_{1}=[-\infty, a), J_{2}=(a, b)$, and $J_{3}=(b, \infty]$ for $a, b \in \mathbb{R}$. By assumption, $\{f \leq b\} \in \mathcal{M}$ and therefore $f^{-1}\left(J_{3}\right)=\{x: f(x) \leq b\}^{\sim} \in \mathcal{M}$. Also,

$$
J_{1}=\bigcup_{k=1}^{\infty}\left[-\infty, a_{k}\right]
$$

where $a_{k}<a$ and $a_{k} \rightarrow a$ as $k \rightarrow \infty$. Hence, $f^{-1}\left(J_{1}\right)=\bigcup_{k=1}^{\infty} f^{-1}\left(\left[-\infty, a_{k}\right]\right)=$ $\bigcup_{k=1}^{\infty}\left\{x: f(x) \leq a_{k}\right\} \in \mathcal{M}$. Finally, $f^{-1}\left(J_{2}\right) \in \mathcal{M}$ since $J_{2}=J_{1} \cap J_{3}$.
5.3. Theorem. A function $f: X \rightarrow \overline{\mathbb{R}}$ is measurable if and only if
(i) $f^{-1}\{-\infty\} \in \mathcal{M}$ and $f^{-1}\{\infty\} \in \mathcal{M}$ and
(ii) $f^{-1}(a, b) \in \mathcal{M}$ for all open intervals $(a, b) \subset \mathbb{R}$.

Proof. If $f$ is measurable, then (i) and (ii) are satisfied since $\{\infty\},\{-\infty\}$ and $(a, b)$ are Borel subsets of $\overline{\mathbb{R}}$.

In order to prove $f$ is measurable, we need to show that $f^{-1}(E) \in \mathcal{M}$ whenever $E$ is a Borel subset of $\overline{\mathbb{R}}$. From (i), and since $E \subset \overline{\mathbb{R}}$ is Borel if and only if $E \cap \mathbb{R}$ is Borel, we only need to show that $f^{-1}(E) \in \mathcal{M}$ whenever $E \subset \mathbb{R}$ is a Borel set. Since $f^{-1}$ preserves unions of sets and since any open set in $\mathbb{R}$ is the disjoint union of open intervals, we see from (ii) that $f^{-1}(U) \in \mathcal{M}$ whenever $U \subset \mathbb{R}$ is an open set. If we define $\Sigma$ as in (5.3) with $Y=\mathbb{R}$, we see that $\Sigma$ is a $\sigma$-algebra that contains the open sets of $\mathbb{R}$ and therefore it contains all Borel sets.

We now proceed to show that measurability is preserved under elementary arithmetic operations on measurable functions. For this, the following will be useful.
5.4. LEMMA. If $f$ and $g$ are measurable functions, then the following sets are measurable:
(i) $X \cap\{x: f(x)>g(x)\}$,
(ii) $X \cap\{x: f(x) \geq g(x)\}$,
(iii) $X \cap\{x: f(x)=g(x)\}$

Proof. If $f(x)>g(x)$, then there is a rational number $r$ such that $f(x)>$ $r>g(x)$. Therefore, it follows that

$$
\{f>g\}=\bigcup_{r \in \mathbb{Q}}(\{f>r\} \cap\{g<r\}),
$$

and (i) easily follows. The set (ii) is the complement of the set (i) with $f$ and $g$ interchanged and it is therefore measurable. The set (iii) is the intersection of two measurable sets of type (ii), and so it too is measurable.

Since all functions under discussion are extended real-valued, we must take some care in defining the sum and product of such functions. If $f$ and $g$ are measurable functions, then $f+g$ is undefined at points where it would be of the form $\infty-\infty$. This difficulty is overcome if we define

$$
(f+g)(x):= \begin{cases}f(x)+g(x) & x \in X-B  \tag{5.5}\\ \alpha & x \in B\end{cases}
$$

where $\alpha \in \overline{\mathbb{R}}$ is chosen arbitrarily and where

$$
\begin{equation*}
B:=\left(f^{-1}\{\infty\} \bigcap g^{-1}\{-\infty\}\right) \bigcup\left(f^{-1}\{-\infty\} \bigcap g^{-1}\{\infty\}\right) \tag{5.6}
\end{equation*}
$$

With this definition we have the following.
5.5. Theorem. If $f, g: X \rightarrow \overline{\mathbb{R}}$ are measurable functions, then $f+g$ and $f g$ are measurable.

Proof. We will treat the case when $f$ and $g$ have values in $\mathbb{R}$. The proof is similar in the general case and is left as an exercise (see Exercise 2, Section 5.1).

To prove that the sum is measurable, define $F: X \rightarrow \mathbb{R} \times \mathbb{R}$ by

$$
F(x)=(f(x), g(x))
$$

and $G: \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}$ by

$$
G(x, y)=x+y
$$

Then $G \circ F(x)=f(x)+g(x)$, so it suffices to show that $G \circ F$ is measurable. Referring to Theorem 5.3, we need only show that $(G \circ F)^{-1}(J) \in \mathcal{M}$ whenever $J \subset \mathbb{R}$ is an open interval. Now $U:=G^{-1}(J)$ is an open set in $\mathbb{R}^{2}$ since $G$ is continuous. Furthermore, $U$ is the union of a countable family, $\mathcal{F}$, of 2-dimensional intervals $I$ of the form $I=I_{1} \times I_{2}$ where $I_{1}$ and $I_{2}$ are open intervals in $\mathbb{R}$. Since

$$
F^{-1}(I)=f^{-1}\left(I_{1}\right) \cap g^{-1}\left(I_{2}\right)
$$

we have

$$
F^{-1}(U)=F^{-1}\left(\bigcup_{I \in \mathcal{F}} I\right)=\bigcup_{I \in \mathcal{F}} F^{-1}(I)
$$

which is a measurable set. Thus $G \circ F$ is measurable since

$$
(G \circ F)^{-1}(J)=F^{-1}(U)
$$

The product is measurable by essentially the same proof.
5.6. Remark. In the situation of abstract measure spaces, if

$$
(X, \mathcal{M}) \xrightarrow{f}(Y, \mathcal{N}) \xrightarrow{g}(Z, \mathcal{P})
$$

are measurable functions, the definitions immediately imply that the composition $g \circ f$ is measurable. Because of this, one might be tempted to conclude that the composition of Lebesgue measurable functions is again Lebesgue measurable. Let's look at this closely. Suppose $f$ and $g$ are Lebesgue measurable functions:

$$
\mathbb{R} \xrightarrow{f} \mathbb{R} \xrightarrow{g} \mathbb{R}
$$

Thus, here we have $X=Y=Z=\mathbb{R}$. Since $Z=\mathbb{R}$, our convention requires that we take $\mathcal{P}$ to be the Borel sets. Moreover, since $f$ is assumed to be Lebesgue measurable, the definition requires $\mathcal{M}$ to be the $\sigma$-algebra of Lebesgue measurable sets. If $g \circ f$ were to be Lebesgue measurable, it would be necessary that $f^{-1}\left(g^{-1}(E)\right)$ is Lebesgue measurable whenever $E$ is a Borel set in $\mathbb{R}$. The definitions imply that it would be necessary $g^{-1}(E)$ to be a Borel set set whenever $E$ is a Borel set in $\mathbb{R}$. The following example shows that this is not generally true.
5.7. Example. (The Cantor-Lebesgue Function) Our example is based on the construction of the Cantor ternary set. Recall (p.94) that the Cantor set $C$ can be expressed as

$$
C=\bigcap_{j=1}^{\infty} C_{j}
$$

where $C_{j}$ is the union of $2^{j}$ closed intervals that remain after the $j^{\text {th }}$ step of the construction. Each of these intervals has length $3^{-j}$. Thus, the set

$$
D_{j}=[0,1]-C_{j}
$$

consists of those $2^{j}-1$ open intervals that are deleted at the $j^{\text {th }}$ step. Let these intervals be denoted by $I_{j, k}, k=1,2, \ldots, 2^{j}-1$, and order them in the obvious way from left to right. Now define a continuous function $f_{j}$ on $[0,1]$ by

$$
\begin{aligned}
& f_{j}(0)=0 \\
& f_{j}(1)=1 \\
& f_{j}(x)=\frac{k}{2^{j}} \quad \text { for } x \in I_{j, k}
\end{aligned}
$$

and define $f_{j}$ linearly on each interval of $C_{j}$. The function $f_{j}$ is continuous, nondecreasing and satisfies

$$
\left|f_{j}(x)-f_{j+1}(x)\right|<\frac{1}{2^{j}} \quad \text { for } x \in[0,1]
$$

Since

$$
\left|f_{j}-f_{j+m}\right|<\sum_{i=j}^{j+m-1} \frac{1}{2^{i}}<\frac{1}{2^{j-1}}
$$

it follows that the sequence $\left\{f_{j}\right\}$ is uniformly Cauchy in the space of continuous functions and thus converges uniformly to a continuous function $f$, called the Cantor-Lebesgue function.

Note that $f$ is nondecreasing and is constant on each interval in the complement of the Cantor set. Furthermore, $f:[0,1] \rightarrow[0,1]$ is onto. In fact, it is easy to see that $f(C)=[0,1]$ because $f(C)$ is compact and $f([0,1]-C)$ is countable.

We use the Cantor-Lebesgue function to show that the composition of Lebesgue measurable functions need not be Lebesgue measurable. Let $h(x)=f(x)+x$ and observe that $h$ is strictly increasing since $f$ is nondecreasing. Thus, $h$ is a homeomorphism from $[0,1]$ onto $[0,2]$. Furthermore, it is clear that $h$ carries the complement of the Cantor set onto an open set of measure 1. Therefore, $h$ maps the Cantor set onto a set $P$ of measure 1. Now let $N$ be a non-Lebesgue measurable subset of $P$; see Exercise 1, Section 4.5. Then, with $A=h^{-1}(N)$, we have $A \subset C$ and therefore $A$ is Lebesgue measurable since $\lambda(A)=0$. Thus we have that $h$ carries a measurable set onto a nonmeasurable set.

Note that $h^{-1}$ is measurable since it is continuous. Let $F:=h^{-1}$. Observe that $A$ is not a Borel set, for if it were, then $F^{-1}(A)$ would be a Borel set. But $F^{-1}(A)=h(A)=N$ and $N$ is not a Borel set. Now $\chi_{A}$ is a Lebesgue measurable function since $A$ is a Lebesgue measurable set. Let $g:=\chi_{A}$. Observe $g^{-1}(1)=A$ and thus $g$ is an example of a Lebesque measurable function that does not preserve Borel sets. Also,

$$
g \circ F=\chi_{A} \circ h^{-1}=\chi_{N}
$$

which shows that this composition of Lebesgue measurable functions is not Lebesgue measurable. To summarize the properties of the Cantor-Lebesgue function, we have:
5.8. Corollary. The Cantor-Lebesgue function, $f$ and its associate $h(x):=$ $f(x)+x$, described above has the following properties:
(i) $f(C)=[0,1]$; that is, $f$ maps a set of measure 0 onto a set of positive measure.
(ii) $h$ maps a Lebesgue measurable set onto a non-measurable set.
(iii) The composition of Lebesgue measurable functions need not be Lebesgue measurable.

Although the example above shows that Lebesgue measurable functions are not generally closed under composition, a positive result can be obtained if the outer
function in the composition is assumed to be Borel measurable. The proof of the following theorem is a direct consequence of the definitions.
5.9. Theorem. Suppose $f: X \rightarrow \overline{\mathbb{R}}$ is measurable and $g: \overline{\mathbb{R}} \rightarrow \overline{\mathbb{R}}$ is Borel measurable. Then $g \circ f$ is measurable. In particular, if $X=\mathbb{R}^{n}$ and $f$ is Lebesgue measurable, then $g \circ f$ is Lebesgue measurable.

The function $g$ is required to have $\overline{\mathbb{R}}$ as its domain of definition because $f$ is extended real-valued; however, any Borel measurable function $g$ defined on $\mathbb{R}$ can be extended to $\overline{\mathbb{R}}$ by assigning arbitrary values to $\infty$ and $-\infty$.

As a consequence of this result, we have the following corollary which complements Theorem 5.5.
5.10. Corollary. Let $f: X \rightarrow \overline{\mathbb{R}}$ be a measurable function.
(i) Let $\varphi(x)=|f(x)|^{p}, 0<p<\infty$, and let $\varphi$ assume arbitrary extended values on the sets $f^{-1}(\infty)$, and $f^{-1}(-\infty)$. Then $\varphi$ is measurable.
(ii) Let $\varphi(x)=\frac{1}{f(x)}$, and let $\varphi$ assume arbitrary extended values on the sets $f^{-1}(0), f^{-1}(\infty)$ and $f^{-1}(-\infty)$. Then $\varphi$ is measurable.

In particular, if $X=\mathbb{R}^{n}$ and $f$ is Lebesgue measurable, then $\varphi$ is Lebesgue measurable in (i) and (ii).

Proof. For (i), define $g(t)=|t|^{p}$ for $t \in \mathbb{R}$ and assign arbitrary values to $g(\infty)$ and $g(-\infty)$. Now apply the previous theorem.

For (ii), proceed in a similar way by defining $g(t)=\frac{1}{t}$ when $t \neq 0, \infty,-\infty$ and assign arbitrary values to $g(0), g(\infty)$, and $g(-\infty)$.

For much of much of the development thus far, the measure $\mu$ in $(X, \mathcal{M}, \mu)$ has played no role. We have only used the fact that $\mathcal{M}$ is a $\sigma$-algebra. Later it will be necessary to deal with functions that are not necessarily defined on all of $X$, but only on the complement of some set of $\mu$-measure 0 . That is, we will deal with functions that are defined only $\mu$-almost everywhere. A measurable set $N$ is called a $\mu$-null set if $\mu(N)=0$. A property that holds for all $x \in X$ except for those $x$ in some $\mu$-null set is said to hold $\mu$-almost everywhere. The term " $\mu$-almost everywhere" is often written in abbreviated form, " $\mu$-a.e.." If it is clear from context that the measure $\mu$ is under consideration, we will simply use the terms "null set" and "almost everywhere."

The next result shows that a measurable function on a complete measure space remains measurable if it is altered on an arbitrary set of measure 0 .
5.11. Theorem. Let $(X, \mathcal{M}, \mu)$ be a complete measure space and let $f, g$ be extended real-valued functions defined on $X$. If $f$ is measurable and $f=g$ almost everywhere, then $g$ is measurable.

Proof. Let $N=\{x: f(x)=g(x)\}$. Then $\mu(\widetilde{N})=0$ and thus, $\widetilde{N}$ as well as all subsets of $\tilde{N}$ are measurable. For $a \in \mathbb{R}$, we have

$$
\begin{aligned}
\{g>a\} & =(\{g>a\} \cap N) \cup(\{g>a\} \cap \tilde{N}) \\
& =(\{f>a\} \cap N) \cup(\{g>a\} \cap \tilde{N}) \in \mathcal{M}
\end{aligned}
$$

5.12. Remark. In case $\mu$ is a complete measure, this result allows us to attach the meaning of measurability to a function $f$ that is defined merely almost everywhere. Indeed, if $N$ is the null set on which $f$ is not defined, we modify the definition of measurability by saying that $f$ is measurable if $\{f>a\} \cap \tilde{N}$ is measurable for each $a \in \mathbb{R}$. This is tantamount to saying that $\bar{f}$ is measurable, where $\bar{f}$ is an extension of $f$ obtained by assigning arbitrary values to $f$ on $N$. This is easily seen because

$$
\{\bar{f}>a\}=(\{\bar{f}>a\} \cap N) \cup(\{f>a\} \cap \tilde{N})
$$

the first set on the right is of measure zero because $\mu$ is complete, and therefore measurable. Furthermore, for functions $f, g$ that are finite-valued at $\mu$-almost every point, we may define $f+g$ as $(f+g)(x)=f(x)+g(x)$ for all those $x \in X$ at which both $f$ and $g$ are defined and do not assume infinite values of opposite sign. Then, if both $f$ and $g$ are measurable, $f+g$ is measurable. A similar discussion holds for the product $f g$.

It therefore becomes apparent that functions that coincide almost everywhere may be considered equivalent. In fact, if we define $f \sim g$ to mean that $f=g$ almost everywhere, then $\sim$ defines an equivalence relation as discussed in Definition 2.9 and thus, a function could be regarded as an equivalence class of functions.

It should be kept in mind that this entire discussion pertains only to the situation in which the measure space $(X, \mathcal{M}, \mu)$ is complete. In particular, it applies in the context of Lebesgue measure on $\mathbb{R}^{n}$, the most important example of a measure space.

We conclude this section by returning to the context of an outer measure $\varphi$ defined on an arbitrary space $X$ as in Definition 4.1. If $f: X \rightarrow \overline{\mathbb{R}}$, then according to Theorem 5.2, $f$ is $\varphi$-measurable if $\{f \leq a\}$ is a $\varphi$-measurable set for each $a \in \mathbb{R}$. That is, with $E_{a}=\{f \leq a\}$, the $\varphi$-measurability of $f$ is equivalent to

$$
\begin{equation*}
\varphi(A)=\varphi\left(A \cap E_{a}\right)+\varphi\left(A-E_{a}\right) \tag{5.7}
\end{equation*}
$$

for any arbitrary set $A \subset X$ and each $a \in \mathbb{R}$. The next result is often useful in applications and gives a characterization of $\varphi$-measurability that appears to be weaker than (5.7).
5.13. Theorem. Suppose $\varphi$ is an outer measure on a space $X$. Then an extended real-valued function $f$ on $X$ is $\varphi$-measurable if and only if

$$
\begin{equation*}
\varphi(A) \geq \varphi(A \cap\{f \leq a\})+\varphi(A \cap\{f \geq b\}) \tag{5.8}
\end{equation*}
$$

whenever $A \subset X$ and $a<b$ are real numbers.

Proof. If $f$ is $\varphi$-measurable, then (5.8) holds since it is implied by (5.7).
To prove the converse, it suffices to show for any real number $r$, that (5.8) implies

$$
E=\{x: f(x) \leq r\}
$$

is $\varphi$-measurable. Let $A \subset X$ be an arbitrary set with $\varphi(A)<\infty$ and define

$$
B_{i}=A \cap\left\{x: r+\frac{1}{i+1} \leq f(x) \leq r+\frac{1}{i}\right\}
$$

for each positive integer $i$. First ${ }^{1}$, we will show

$$
\begin{equation*}
\infty>\varphi(A) \geq \varphi\left(\bigcup_{k=1}^{\infty} B_{2 k}\right)=\sum_{k=1}^{\infty} \varphi\left(B_{2 k}\right) \tag{5.9}
\end{equation*}
$$

The proof is by induction, so assume (5.9) is valid as $k$ runs from 1 to $j-1$. That is, assume

$$
\begin{equation*}
\varphi\left(\bigcup_{k=1}^{j-1} B_{2 k}\right)=\sum_{k=1}^{j-1} \varphi\left(B_{2 k}\right) \tag{5.10}
\end{equation*}
$$

Let

$$
A_{j}=\bigcup_{k=1}^{j-1} B_{2 k}
$$

Then, using (5.8), the induction hypothesis, and the fact that

$$
\begin{equation*}
\left(B_{2 j} \cup A_{j}\right) \cap\left\{f \leq r+\frac{1}{2 j}\right\}=B_{2 j} \tag{5.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(B_{2 j} \cup A_{j}\right) \cap\left\{f \geq r+\frac{1}{2 j-1}\right\}=A_{j} \tag{5.12}
\end{equation*}
$$

[^0]we obtain
\[

$$
\begin{align*}
& \varphi\left(\bigcup_{k=1}^{j} B_{2 k}\right)=\varphi\left(B_{2 j} \cup A_{j}\right) \\
& \quad \begin{array}{ll} 
& \geq \varphi\left[\left(B_{2 j} \cup A_{j}\right) \cap\left\{f \leq r+\frac{1}{2 j}\right\}\right] \\
& +\varphi\left[\left(B_{2 j} \cup A_{j}\right) \cap\left\{f \geq r+\frac{1}{2 j-1}\right\}\right] \\
& =\varphi\left(B_{2 j}\right)+\varphi\left(A_{j}\right) \\
& =\varphi\left(B_{2 j}\right)+\sum_{k=1}^{j-1} \varphi\left(B_{2 k}\right) \quad \text { by }(5.11) \text { and }(5.12) \\
\quad=\sum_{k=1}^{j} \varphi\left(B_{2 k}\right) .
\end{array} \quad \text { by induction hypothesis }(5.10)
\end{align*}
$$
\]

Thus, (5.9) is valid as $k$ runs from 1 to $j$ for any positive integer $j$. In other words, we obtain

$$
\infty>\varphi(A) \geq \varphi\left(\bigcup_{k=1}^{\infty} B_{2 k}\right) \geq \varphi\left(\bigcup_{k=1}^{j} B_{2 k}\right)=\sum_{k=1}^{j} \varphi\left(B_{2 k}\right)
$$

for any positive integer $j$. This implies

$$
\infty>\varphi(A) \geq \sum_{k=1}^{\infty} \varphi\left(B_{2 k}\right)
$$

Virtually the same argument can be used to obtain

$$
\infty>\varphi(A) \geq \sum_{k=1}^{\infty} \varphi\left(B_{2 k-1}\right)
$$

thus implying

$$
\infty>2 \varphi(A) \geq \sum_{k=1}^{\infty} \varphi\left(B_{k}\right)
$$

Now the tail end of this convergent series can be made arbitrarily small; that is, for each $\varepsilon>0$ there exists a positive integer $m$ such that

$$
\begin{aligned}
\varepsilon>\sum_{i=m}^{\infty} \varphi\left(B_{i}\right) & =\varphi\left(\bigcup_{i=m}^{\infty} B_{i}\right) \\
& \geq \varphi\left(A \cap\left\{r<f<r+\frac{1}{m}\right\}\right)
\end{aligned}
$$

For ease of notation, define an outer measure $\psi(S)=\varphi(S \cap A)$ whenever $S \subset X$.
With this notation, we have shown

$$
\begin{aligned}
\varepsilon & >\psi\left(\left\{r<f<r+\frac{1}{m}\right\}\right) \\
& =\psi\left(\{r<f\} \cap\left\{f<r+\frac{1}{m}\right\}\right) \\
& \geq \psi(\{f>r\})-\psi\left(\left\{f \geq r+\frac{1}{m}\right\}\right) .
\end{aligned}
$$

The last inequality is implied by the subadditivity of $\psi$. Therefore,

$$
\begin{align*}
\varphi(A \cap E)+\varphi(A-E) & =\psi(E)+\psi(\widetilde{E}) \\
& =\psi(E)+\psi(\{f>r\}) \\
& \leq \psi(E)+\psi\left(\left\{f \geq r+\frac{1}{m}\right\}\right)+\varepsilon \\
& =\varphi(A \cap E)+\varphi\left(A \cap\left\{f \geq r+\frac{1}{m}\right\}\right)+\varepsilon \\
& \leq \varphi(A)+\varepsilon \tag{5.8}
\end{align*}
$$

Since $\varepsilon$ is arbitrary, this proves that $E$ is $\varphi$-measurable.

## Exercises for Section 5.1

1. Let $(X, \mathcal{M}) \xrightarrow{f}(Y, \mathcal{B})$ be a continuous mapping where $X$ and $Y$ are topological spaces, $\mathcal{M}$ is a $\sigma$-algebra that contains the Borel sets in $X$ and $\mathcal{B}$ is the family of Borel sets in $Y$. Prove that $f$ is measurable.
2. Complete the proof of Theorem 5.5 when $f$ and $g$ have values in $\overline{\mathbb{R}}$.
3. Prove that a function defined on $\mathbb{R}^{n}$ that is continuous everywhere except for a set of Lebesgue measure zero is a Lebesgue measurable function. In particular, conclude that a nondecreasing function defined on $[0,1]$ is Lebesgue measurable.

### 5.2. Limits of Measurable Functions


#### Abstract

In order to be useful in applications, it is necessary for measurability to be preserved by virtually all types of limit operations on sequences of measurable functions. In this section, it is shown that measurability is preserved under the operations of upper and lower limits of sequences of functions as well as upper and lower envelopes. It is also shown that on a finite measure space, pointwise a.e. convergence of a sequence of measurable functions implies uniform convergence on the complements of sets of arbitrarily small measure (Egoroff's Theorem). Finally, the relationship between convergence in measure and pointwise a.e. convergence is investigated.


Throughout this section, it will be assumed that all functions are $\overline{\mathbb{R}}$-valued, unless otherwise stated.
5.14. Definition. Let $(X, \mathcal{M}, \mu)$ be a measure space, and let $\left\{f_{i}\right\}$ be a sequence of measurable functions defined on $X$. The upper and lower envelopes of $\left\{f_{i}\right\}$ are defined respectively as

$$
\sup _{i} f_{i}(x)=\sup \left\{f_{i}(x): i=1,2, \ldots\right\}
$$

and

$$
\inf _{i} f_{i}(x)=\inf \left\{f_{i}(x): i=1,2, \ldots\right\} .
$$

Also, the upper and lower limits of $\left\{f_{i}\right\}$ are defined as

$$
\limsup _{i \rightarrow \infty} f_{i}(x)=\inf _{j \geq 1}\left(\sup _{i \geq j} f_{i}(x)\right)
$$

and

$$
\liminf _{i \rightarrow \infty} f_{i}(x)=\sup _{j \geq 1}\left(\inf _{i \geq j} f_{i}(x)\right) .
$$

5.15. Theorem. Let $\left\{f_{i}\right\}$ be a sequence of measurable functions defined on the measure space $(X, \mathcal{M}, \mu)$. Then $\sup _{i} f_{i}, \inf _{i} f_{i}, \limsup _{i \rightarrow \infty} f_{i}$, and $\liminf _{i \rightarrow \infty} f_{i}$ are all measurable functions.

Proof. For each $a \in \mathbb{R}$ the identity

$$
X \cap\left\{x: \sup _{i} f_{i}(x)>a\right\}=\bigcup_{i=1}^{\infty}\left(X \cap\left\{f_{i}(x)>a\right\}\right)
$$

implies that $\sup f_{i}$ is measurable. The measurability of the lower envelope follows from

$$
\inf _{i} f_{i}(x)=-\sup _{i}\left(-f_{i}(x)\right) .
$$

Now that it has been shown that the upper and lower envelopes are measurable, it is immediate that the upper and lower limits of $\left\{f_{i}\right\}$ are also measurable.

We begin by investigating what information can be deduced from the pointwise almost everywhere convergence of a sequence of measurable functions on a finite measure space.
5.16. Definition. A sequence of measurable functions, $\left\{f_{i}\right\}$, with the property that

$$
\lim _{i \rightarrow \infty} f_{i}(x)=f(x)
$$

for $\mu$-almost every $x \in X$ is said to converge pointwise almost everywhere (or more briefly, converge pointwise a.e.) to $f$.

We have the following:
5.17. Corollary. Let $X=\mathbb{R}^{n}$. If $\left\{f_{i}\right\}$ is a sequence of Lebesgue measurable functions that converge pointwise almost everywhere to $f$, then $f$ is measurable.

The following is one of the main results of this section.
5.18. Theorem (Egoroff). Let $(X, \mathcal{M}, \mu)$ be a finite measure space and suppose $\left\{f_{i}\right\}$ and $f$ are measurable functions that are finite almost everywhere on $X$. Also, suppose that $\left\{f_{i}\right\}$ converges pointwise a.e. to $f$. Then for each $\varepsilon>0$ there exists a set $A \in \mathcal{M}$ such that $\mu(\widetilde{A})<\varepsilon$ and $\left\{f_{i}\right\} \rightarrow f$ uniformly on $A$.

First, we will prove the following lemma.
5.19. ThEOREM (Egoroff). Assume the hypotheses of the previous theorem. Then for each pair of numbers $\varepsilon, \delta>0$, there exist a set $A \in \mathcal{M}$ and an integer $i_{0}$ such that $\mu(\widetilde{A})<\varepsilon$ and

$$
\left|f_{i}(x)-f(x)\right|<\delta
$$

whenever $x \in A$ and $i \geq i_{0}$.
Proof. Choose $\varepsilon, \delta>0$. Let $E$ denote the set on which the functions $f_{i}, i=$ $1,2, \ldots$, and $f$ are defined and finite. Also, let $F$ be the set on which $\left\{f_{i}\right\}$ converges pointwise to $f$. With $A_{0}:=E \cap F$, we have by hypothesis, $\mu\left(\widetilde{A}_{0}\right)=0$. For each positive integer $i$, let

$$
A_{i}=A_{0} \cap\left\{x:\left|f_{j}(x)-f(x)\right|<\delta \text { for all } j \geq i\right\}
$$

Then, $A_{1} \subset A_{2} \subset \ldots$ and $\cup_{i=1}^{\infty} A_{i}=A_{0}$ and consequently, $\widetilde{A}_{1} \supset \widetilde{A}_{2} \supset \ldots$ with $\cap_{i=1}^{\infty} \widetilde{A}_{i}=\widetilde{A}_{0}$. Since $\mu\left(\widetilde{A}_{1}\right) \leq \mu(X)<\infty$, it follows from Theorem 4.49 (v) that

$$
\lim _{i \rightarrow \infty} \mu\left(\widetilde{A}_{i}\right)=\mu\left(\widetilde{A}_{0}\right)=0
$$

The result follows by choosing $i_{0}$ such that $\mu\left(\widetilde{A}_{i_{0}}\right)<\varepsilon$ and $A=A_{i_{0}}$.
Proof of Egoroff's Theorem. Choose $\varepsilon>0$. By the previous lemma, for each positive integer $i$, there exist a positive integer $j_{i}$ and a measurable set $A_{i}$ such that

$$
\mu\left(\widetilde{A}_{i}\right)<\frac{\varepsilon}{2^{i}} \quad \text { and } \quad\left|f_{j}(x)-f(x)\right|<\frac{1}{i}
$$

for all $x \in A_{i}$ and all $j \geq j_{i}$. With $A$ defined as $A=\cap_{i=1}^{\infty} A_{i}$, we have

$$
\widetilde{A}=\bigcup_{i=1}^{\infty} \widetilde{A}_{i}
$$

and

$$
\mu(\widetilde{A}) \leq \sum_{i=1}^{\infty} \mu\left(\widetilde{A}_{i}\right)<\sum_{i=1}^{\infty} \frac{\varepsilon}{2^{i}}=\varepsilon
$$

Furthermore, if $j \geq j_{i}$, then

$$
\sup _{x \in A}\left|f_{j}(x)-f(x)\right| \leq \sup _{x \in A_{i}}\left|f_{j}(x)-f(x)\right| \leq \frac{1}{i}
$$

for every positive integer $i$. This implies that $\left\{f_{i}\right\} \rightarrow f$ uniformly on $A$.
5.20. Corollary. In the previous theorem, assume in addition that $X$ is a metric space and that $\mu$ is a Borel measure with $\mu(X)<\infty$. Then A can be taken as a closed set.

Proof. The previous theorem provides a set $A$ such that $A \in \mathcal{M},\left\{f_{i}\right\}$ converges uniformly to $f$ and $\mu(\tilde{A})<\varepsilon / 2$. Since $\mu$ is a finite Borel measure, we see from Theorem 4.63 (p.126) that there exists a closed set $F \subset A$ with $\mu(A \backslash F)<\varepsilon / 2$. Hence, $\mu(\tilde{F})<\varepsilon$ and $\left\{f_{i}\right\} \rightarrow f$ uniformly on $F$.
5.21. Definition. Because of its importance, we attach a name to the type of convergence exhibited in the conclusion of Egoroff's Theorem. Suppose that $\left\{f_{i}\right\}$ and $f$ are measurable functions that are finite almost everywhere. We say that $\left\{f_{i}\right\}$ converges to $f$ almost uniformly if for every $\varepsilon>0$, there exists a set $A \in \mathcal{M}$ such that $\mu(\widetilde{A})<\varepsilon$ and $\left\{f_{i}\right\}$ converges to $f$ uniformly on $A$. Thus, Egoroff's Theorem states that pointwise a.e. convergence on a finite measure space implies almost uniform convergence. The converse is also true and is left as Exercise 5, Section 5.2.
5.22. Remark. The hypothesis that $\mu(X)<\infty$ is essential in Egoroff's Theorem. Consider the case of Lebesgue measure on $\mathbb{R}$ and define a sequence of functions by

$$
f_{i}=\chi_{[i, \infty)}
$$

for each positive integer $i$. Then, $\lim _{i \rightarrow \infty} f_{i}(x)=0$ for each $x \in \mathbb{R}$, but $\left\{f_{i}\right\}$ does not converge uniformly to 0 on any set $A$ whose complement has finite Lebesgue measure. Indeed, for any such set, it would follow that $\widetilde{A}$ does not contain any $[i, \infty)$; that is, for each $i$, there would exist $x \in[i, \infty) \cap A$ with $f_{i}(x)=1$, thus showing that $\left\{f_{i}\right\}$ does not converge uniformly to 0 on $A$.
5.23. Definition. A sequence of measurable functions $\left\{f_{i}\right\}$ defined relative to the measure space $(X, \mathcal{M}, \mu)$ is said to converge in measure to a measurable function $f$ if for every $\varepsilon>0$, we have

$$
\lim _{i \rightarrow \infty} \mu\left(X \cap\left\{x:\left|f_{i}(x)-f(x)\right| \geq \varepsilon\right\}\right)=0
$$

We already encountered a result (Lemma 5.19) that essentially shows that pointwise a.e. convergence on a finite measure space implies convergence in measure. Formally, it is as follows.
5.24. Theorem. Let $(X, \mathcal{M}, \mu)$ be a finite measure space, and suppose $\left\{f_{i}\right\}$ and $f$ are measurable functions that are finite a.e. on $X$. If $\left\{f_{i}\right\}$ converges to $f$ a.e. on $X$, then $\left\{f_{i}\right\}$ converges to $f$ in measure.

Proof. Choose positive numbers $\varepsilon$ and $\delta$. According to Lemma 5.19, there exist a set $A \in \mathcal{M}$ and an integer $i_{0}$ such that $\mu(\widetilde{A})<\varepsilon$ and

$$
\left|f_{i}(x)-f(x)\right|<\delta
$$

whenever $x \in A$ and $i \geq i_{0}$. Thus,

$$
X \cap\left\{x:\left|f_{i}(x)-f(x)\right| \geq \delta\right\} \subset \widetilde{A}
$$

if $i \geq i_{0}$. Since $\mu(\widetilde{A})<\varepsilon$ and $\varepsilon>0$ is arbitrary, the result follows.
5.25. Remark. It is easy to see that the converse is not true. Let $X=[0,1]$ with $\mu$ taken as Lebesgue measure. Consider a sequence of partitions of $[0,1], \mathcal{P}_{i}$, each consisting of closed, nonoverlapping intervals of length $1 / 2^{i}$. Let $\mathcal{F}$ denote the family of all intervals comprising the partitions $\mathcal{P}_{i}, i=1,2, \ldots$. Linearly order $\mathcal{F}$ by defining $I \leq I^{\prime}$ if both $I$ and $I^{\prime}$ are elements of the same partition $\mathcal{P}_{i}$ and if $I$ is to the left of $I^{\prime}$. Otherwise, define $I \leq I^{\prime}$ if the length of $I$ is no greater than that of $I^{\prime}$. Now put the elements of $\mathcal{F}$ into a one-to-one order preserving correspondence with the positive integers. With the elements of $\mathcal{F}$ labeled as $I_{k}, k=1,2, \ldots$, define a sequence of functions $\left\{f_{k}\right\}$ by $f_{k}=\chi_{I_{k}}$. Then it is easy to see that $\left\{f_{k}\right\} \rightarrow 0$ in measure but that $\left\{f_{k}(x)\right\}$ does not converge to 0 for any $x \in[0,1]$.

Although the sequence $\left\{f_{k}\right\}$ converges nowhere to 0 , it does have a subsequence that converges to 0 a.e., namely the subsequence

$$
f_{1}, f_{2}, f_{4}, \ldots, f_{2^{k-1}}, \ldots
$$

In fact, this sequence converges to 0 at all points except $x=0$. This illustrates the following general result.
5.26. Theorem. Let $(X, \mathcal{M}, \mu)$ be a measure space and let $\left\{f_{i}\right\}$ and $f$ be measurable functions such that $f_{i} \rightarrow f$ in measure. Then there exists a subsequence $\left\{f_{i_{j}}\right\}$ such that

$$
\lim _{j \rightarrow \infty} f_{i_{j}}(x)=f(x)
$$

for $\mu$-a.e. $x \in X$
Proof. Let $i_{1}$ be a positive integer such that

$$
\mu\left(X \cap\left\{x:\left|f_{i_{1}}(x)-f(x)\right| \geq 1\right\}\right)<\frac{1}{2}
$$

Assuming that $i_{1}, i_{2}, \ldots, i_{k}$ have been chosen, let $i_{k+1}>i_{k}$ be such that

$$
\mu\left(X \cap\left\{x:\left|f_{i_{k+1}}(x)-f(x)\right| \geq \frac{1}{k+1}\right\}\right) \leq \frac{1}{2^{k+1}}
$$

Let

$$
A_{j}=\bigcup_{k=j}^{\infty}\left\{x:\left|f_{i_{k}}(x)-f(x)\right| \geq \frac{1}{k}\right\}
$$

and observe that the sequence $A_{j}$ is descending. Since

$$
\mu\left(A_{1}\right)<\sum_{k=1}^{\infty} \frac{1}{2^{k}}<\infty
$$

with $B=\cap_{j=1}^{\infty} A_{j}$, it follows that

$$
\mu(B)=\lim _{j \rightarrow \infty} \mu\left(A_{j}\right) \leq \lim _{j \rightarrow \infty} \sum_{k=j}^{\infty} \frac{1}{2^{k}}=\lim _{j \rightarrow \infty} \frac{1}{2^{j-1}}=0 .
$$

Now select $x \in \widetilde{B}$. Then there exists an integer $j=j_{x}$ such that

$$
x \in \widetilde{A_{j_{x}}}=\bigcap_{k=j_{x}}^{\infty}\left(X \cap\left\{y:\left|f_{i_{k}}(y)-f(y)\right|<\frac{1}{k}\right\}\right)
$$

If $\varepsilon>0$, choose $k_{0}$ such that $k_{0} \geq j_{x}$ and $\frac{1}{k_{0}} \leq \varepsilon$. Then for $k \geq k_{0}$, we have

$$
\left|f_{i_{k}}(x)-f(x)\right|<\frac{1}{k} \leq \varepsilon
$$

which implies that $f_{i_{k}}(x) \rightarrow f(x)$ for all $x \in \widetilde{B}$.
There is another mode of convergence, fundamental in measure, which is discussed in Exercise 6, Section 5.2.

## Exercises for Section 5.2

1. Let $\mathcal{F}$ be a family of continuous functions on a metric space $(X, \rho)$. Let $f$ denote the upper envelope of the family $\mathcal{F}$; that is,

$$
f(x)=\sup \{g(x): g \in \mathcal{F}\} .
$$

Prove for each real number $a$, that $\{x: f(x)>a\}$ is open.
2. Let $f(x, y)$ be a function defined on $\mathbb{R}^{2}$ that is continuous in each variable separately. Prove that $f$ is Lebesgue measurable. Hint: Approximate $f$ in the variable $x$ by piecewise-linear continuous functions $f_{n}$ so that $f_{n} \rightarrow f$ pointwise.
3. Let $(X, \mathcal{M}, \mu)$ be a finite measure space. Suppose that $\left\{f_{i}\right\}_{i=1}^{\infty}$ and $f$ are measurable functions. Prove that $f_{i} \rightarrow f$ in measure if and only if each subsequence of $f_{i}$ has a subsequence that converges to $f \mu$-a.e.
4. Show that the supremum of an uncountable family of measurable $\overline{\mathbb{R}}$-valued functions can fail to be measurable.
5. Suppose $(X, \mathcal{M}, \mu)$ is a finite measure space. Prove that almost uniform convergence implies convergence almost everywhere.
6. A sequence $\left\{f_{i}\right\}$ of a.e. finite-valued measurable functions on a measure space $(X, \mathcal{M}, \mu)$ is fundamental in measure if, for every $\varepsilon>0$,

$$
\mu\left(\left\{x:\left|f_{i}(x)-f_{j}(x)\right| \geq \varepsilon\right\}\right) \rightarrow 0
$$

as $i$ and $j \rightarrow \infty$. Prove that if $\left\{f_{i}\right\}$ is fundamental in measure, then there is a measurable function $f$ to which the sequence $\left\{f_{i}\right\}$ converges in measure. Hint: Choose integers $i_{j+1}>i_{j}$ such that $\mu\left\{\left|f_{i_{j}}-f_{i_{j+1}}\right|>2^{-j}\right\}<2^{-j}$. The sequence $\left\{f_{i_{j}}\right\}$ converges a.e. to a function $f$. Then it follows that

$$
\left\{\left|f_{i}-f\right| \geq \varepsilon\right\} \subset\left\{\left|f_{i}-f_{i_{j}}\right| \geq \varepsilon / 2\right\} \cup\left\{\left|f_{i_{j}}-f\right| \geq \varepsilon / 2\right\}
$$

By hypothesis, the measure of the first term on the right is arbitrarily small if $i$ and $i_{j}$ are large, and the measure of the second term tends to 0 since almost uniform convergence implies convergence in measure.

### 5.3. Approximation of Measurable Functions

In Section 3.2 certain fundamental approximation properties of Carathéodory outer measures were established. In particular, it was shown that each Borel set $B$ of finite measure contains a closed set whose measure is arbitrarily close to that of $B$. The structure of a Borel set can be very complicated, but yet this result states that a complicated set can be approximated by one with an elementary topological property. In this section we pursue an analogous situation by showing that each measurable function on a metric space of finite measure is almost continuous (Lusin's Theorem). That is, every measurable function is continuous on sets whose complements have arbitrarily small measure. This result is in the same spirit as Egoroff's Theorem, which states that pointwise a.e. convergence implies almost uniform convergence on a finite measure space.

The characteristic function of a measurable set is the most elementary example of a measurable function. The next level of complexity involves linear combinations of such functions. A simple function on $X$ is one that assumes only a finite number of values: Thus, the range of a simple function, $f$, is a finite subset of $\mathbb{R}$. If rng $f=\left\{a_{1}, a_{2}, \ldots, a_{k}\right\}$, and $A_{i}=f^{-1}\left\{a_{i}\right\}$, then $f$ can be written as

$$
f=\sum_{i=1}^{k} a_{i} \chi_{A_{i}} .
$$

If $X=\mathbb{R}^{n}$, a step function is of the form $f=\sum_{k=1}^{N} a_{k} \chi_{R_{k}}$, where each $R_{k}$ is an interval and the $a_{k}$ are real numbers.

We begin by proving that any measurable function is the pointwise a.e. limit of measurable simple functions.
5.27. ThEOREM. Let $f: X \rightarrow \overline{\mathbb{R}}$ be an arbitrary (possibly nonmeasurable) function. Then the following hold.
(i) There exists a sequence of simple functions, $\left\{f_{i}\right\}$, such that $f_{i}(x) \rightarrow f(x)$ for each $x \in X$,
(ii) If $f$ is nonnegative, the sequence can be chosen so that $f_{i} \uparrow f$,
(iii) If $f$ is bounded, the sequence can be chosen so that $f_{i} \rightarrow f$ uniformly on $X$,
(iv) If $f$ is measurable, the $f_{i}$ can be chosen to be measurable.

Proof. Assume first that $f \geq 0$. For each positive integer $i$, partition $[0, i)$ into $i \cdot 2^{i}$ half-open intervals of the form $\left[\frac{k-1}{2^{i}}, \frac{k}{2^{i}}\right), k=1,2, \ldots, i \cdot 2^{i}$. Label these intervals as $H_{i, k}$ and let

$$
A_{i, k}=f^{-1}\left(H_{i, k}\right) \quad \text { and } \quad A_{i}=f^{-1}([i, \infty])
$$

These sets are pairwise disjoint and form a partition of $X$. The approximating simple function $f_{i}$ on $X$ is defined as

$$
f_{i}(x)=\left\{\begin{array}{cl}
\frac{k-1}{2^{i}} & x \in A_{i, k} \\
i & x \in A_{i}
\end{array}\right.
$$

If $f$ is measurable, then the sets $A_{i, k}$ and $A_{i}$ are measurable and thus, so are the functions $f_{i}$. Moreover, it is easy to see that

$$
f_{1} \leq f_{2} \leq \ldots \leq f
$$

If $f(x)<\infty$, then for every $i>f(x)$ we have

$$
\left|f_{i}(x)-f(x)\right|<\frac{1}{2^{i}}
$$

and hence $f_{i}(x) \rightarrow f(x)$. If $f(x)=\infty$ then $f_{i}(x)=i \rightarrow f(x)$. In any case we obtain

$$
\lim _{i \rightarrow \infty} f_{i}(x)=f(x) \text { for } x \in X
$$

Suppose $f$ is bounded by some number, say $M$; that is, suppose $f(x) \leq M$ for all $x \in X$. Then $A_{i}=\emptyset$ for all $i>M$ and therefore $\left|f_{i}(x)-f(x)\right|<1 / 2^{i}$ for all $x \in X$, thus showing that $f_{i} \rightarrow f$ uniformly if $f$ is bounded. This establishes the Theorem in case $f \geq 0$.

In general, let $f^{+}(x)=\max (f(x), 0)$ and $f^{-}(x)=-\min (f(x), 0)$ denote the positive and negative parts of $f$. Then $f^{+}$and $f^{-}$are nonnegative and $f=f^{+}-$ $f^{-}$. Now apply the previous results to $f^{+}$and $f^{-}$to obtain the final form of the Theorem.

The proof of the following Corollary is left to the reader (see Exercise 2, Section 5.3).
5.28. Corollary. Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a Lebesgue measurable function. Then, there exists a sequence of step functions $f_{i}$ such that $f_{i}(x) \rightarrow f(x)$ for almost every $x \in \mathbb{R}^{n}$.

Since it is possible for a measurable function to be discontinuous at every point of its domain, it seems unlikely that an arbitrary measurable function would have any regularity properties. However, the next result gives some information in the positive direction. It states, roughly, that any measurable function $f$ is continuous on a closed set $F$ whose complement has arbitrarily small measure. It is important to note that the result asserts the function is continuous on $F$ with respect to the relative topology on $F$. It should not be interpreted to say that $f$ is continuous at every point of $F$ relative to the topology on $X$.
5.29. Theorem (Lusin's Theorem). Suppose $(X, \mathcal{M}, \mu)$ is a measure space where $X$ is a metric space and $\mu$ is a finite Borel measure. Let $f: X \rightarrow \overline{\mathbb{R}}$ be a measurable function that is finite almost everywhere. Then for every $\varepsilon>0$ there is a closed set $F \subset X$ with $\mu(\widetilde{F})<\varepsilon$ such that $f$ is continuous on $F$ in the relative topology.

Proof. Choose $\varepsilon>0$. For each fixed positive integer $i$, write $\mathbb{R}$ as the disjoint union of half-open intervals $H_{i, j}, j=1,2, \ldots$, whose lengths are $1 / i$. Consider the disjoint measurable sets

$$
A_{i, j}=f^{-1}\left(H_{i, j}\right)
$$

and refer to Theorem 4.63 to obtain disjoint closed sets $F_{i, j} \subset A_{i, j}$ such that $\mu\left(A_{i, j}-F_{i, j}\right)<\varepsilon / 2^{i+j}, j=1,2, \ldots$. Let

$$
E_{k}=X-\bigcup_{j=1}^{k} F_{i, j}
$$

for $k=1,2, \ldots, \infty$. (Keep in mind that $i$ is fixed, so it is not necessary to indicate that $E_{k}$ depends on $i$ ). Then $E_{1} \supset E_{2} \supset \ldots, \cap_{k=1}^{\infty} E_{k}=E_{\infty}$, and

$$
\mu\left(E_{\infty}\right)=\mu\left(X-\bigcup_{j=1}^{\infty} F_{i, j}\right)=\sum_{j=1}^{\infty} \mu\left(A_{i, j}-F_{i, j}\right)<\frac{\varepsilon}{2^{i}}
$$

Since $\mu(X)<\infty$, it follows that

$$
\lim _{k \rightarrow \infty} \mu\left(E_{k}\right)=\mu\left(E_{\infty}\right)<\frac{\varepsilon}{2^{i}}
$$

Hence, there exists a positive integer $J=J(i)$ such that

$$
\mu\left(E_{J}\right)=\mu\left(X-\bigcup_{j=1}^{J} F_{i, j}\right)<\frac{\varepsilon}{2^{i}}
$$

For each $H_{i, j}$, select an arbitrary point $y_{i, j} \in H_{i, j}$ and let $B_{i}=\cup_{j=1}^{J} F_{i, j}$. Then define a continuous function $g_{i}$ on the closed set $B_{i}$ by

$$
g_{i}(x)=y_{i, j} \quad \text { whenever } \quad x \in F_{i, j}, j=1,2, \ldots, J
$$

The functions $g_{i}$ are continuous (relative to $B_{i}$ ) because the closed sets $F_{i, j}$ are disjoint. Note that $\left|f(x)-g_{i}(x)\right|<1 / i$ for $x \in B_{i}$. Therefore, on the closed set

$$
F=\bigcap_{i=1}^{\infty} B_{i} \quad \text { with } \quad \mu(X-F) \leq \sum_{i=1}^{\infty} \mu\left(X-B_{i}\right)<\varepsilon,
$$

it follows that the continuous functions $g_{i}$ converge uniformly to $f$, thus proving that $f$ is continuous on $F$.

Using Corollary 4.56 we can rewrite Lusin's Theorem as follows:
5.30. Corollary. Let $\varphi$ be a Borel regular outer measure on a metric space $X$. Let $\mathcal{M}$ be the $\sigma$-algebra of $\varphi$-measurable sets. Consider the measure space $(X, \mathcal{M}, \varphi)$ and let $A \in \mathcal{M}, \varphi(A)<\infty$. If $f: A \rightarrow \overline{\mathbb{R}}$ is a measurable function that is finite almost everywhere, then for every $\epsilon>0$ there exists a closed set $F \subset A$ with $\varphi(A \backslash F)<\epsilon$ such that $f$ is continuous on $F$ in the relative topology.

In particular, since $\lambda^{*}$ is a Borel regular outer measure, we conclude that if $f: A \rightarrow \mathbb{R}, A \subset \mathbb{R}^{n}$ Lebesgue measurable, is a Lebesgue measurable function with $\lambda(A)<\infty$, then for every $\epsilon>0$ there exists a closed set $F \subset A$ with $\lambda(A \backslash F)<\epsilon$ such that $f$ is continuous on $F$ in the relative topology.

We close this chapter with a table that reflects the interaction of the various types of convergences that we have encountered so far. A convergence-type listed in the first column implies one in the first row if the corresponding entry of the matrix is indicated by $\uparrow$ (along with the appropriate hypothesis).

|  | Fundamental in measure | Convergence in measure | Almost uniform convergence | Pointwise a．e．conver－ gence |
| :---: | :---: | :---: | :---: | :---: |
| Fundamental in measure | 介 | $\Uparrow$ | 介 <br> For a sub－ sequence if $\mu(X)<\infty$ | $\Uparrow$ <br> For a subse－ quence |
| Convergence in measure | 介 | $\Uparrow$ | $\Uparrow$ For a sub－ sequence if $\mu(X)<\infty$ | $\Uparrow$ <br> For a subse－ quence |
| Almost uniform convergence | 介 | $\Uparrow$ | 介 | 介 |
| Pointwise a．e．conver－ gence | if $\mu(X)<\infty$ | 介 if $\mu(X)<\infty$ | 介 if $\mu(X)<\infty$ | 介 |

## Exercises for Section 5.3

1．Let $E \subset \mathbb{R}^{n}$ be a Lebesgue measurable set with $\lambda(E)<\infty$ and let $\chi_{E}$ be the characteristic function of $E$ ．Prove that there is a sequence of step functions $\left\{\psi_{k}\right\}_{k=1}^{\infty}$ that converges pointwise to $\chi_{E}$ almost everywhere．Hint：Show that if $\lambda(E)<\infty$ ，then there exists a finite union of closed intervals $Q_{j}$ such that $F=\cup_{j=1}^{N} Q_{j}$ and $\lambda(E \Delta F) \leq \epsilon$ ．Recall that $E \Delta F=(E \backslash F) \cup(F \backslash E)$ ．
2．Use the previous exercise 1 to prove Corollary 5．28．
3．A union of $n$－dimensional（closed）intervals in $\mathbb{R}^{n}$ is said to be almost disjoint if the interiors of the intervals are disjoint．Show that every open subset $U$ of $\mathbb{R}^{n}, n \geq 1$ ，can be written as a countable union of almost disjoint intervals．
4．Let $(X, \mathcal{M}, \mu)$ be a $\sigma$－finite measure space and suppose that $f, f_{k}, k=1,2, \ldots$ ， are measurable functions that are finite almost everywhere and

$$
\lim _{k \rightarrow \infty} f_{k}(x)=f(x)
$$

for $\mu$ almost all $x \in X$ ．Prove that there are measurable sets
$E_{0}, E_{1}, E_{2}, \ldots$, such that $\nu\left(E_{0}\right)=0$,

$$
X=\bigcup_{i=0}^{\infty} E_{i},
$$

and $\left\{f_{k}\right\} \rightarrow f$ uniformly on each $E_{i}, i>0$.
5. Use Corollary 4.56 to prove Corollary 5.30.
6. Suppose $f:[0,1] \rightarrow \mathbb{R}$ is Lebesgue measurable. For any $\varepsilon>0$, show that there is a continuous function $g$ on $[0,1]$ such that

$$
\lambda([0,1] \cap\{x: f(x) \neq g(x)\})<\varepsilon .
$$

## CHAPTER 6

## Integration

### 6.1. Definitions and Elementary Properties


#### Abstract

Based on the ideas of H. Lebesgue, a far-reaching generalization of Riemann integration has been developed. In this section we define and deduce the elementary properties of integration with respect to an abstract measure.


We first extend the notion of simple function to allow better approximation of unbounded functions. Throughout this section and the next, we will assume the context of a general measure space $(X, \mathcal{M}, \mu)$.
6.1. Definition. A function $f: X \rightarrow \overline{\mathbb{R}}$ is called countably-simple if it assumes only a countable number of values, including possibly $\pm \infty$. Given a measure space $(X, \mathcal{M}, \mu)$, the integral of a nonnegative measurable countablysimple function $f: X \rightarrow \overline{\mathbb{R}}$ is defined to be

$$
\int_{X} f d \mu=\sum_{i=1}^{\infty} a_{i} \mu\left(f^{-1}\left\{a_{i}\right\}\right)
$$

where the range of $f=\left\{a_{1}, a_{2}, \ldots\right\}$ and, by convention, $0 \cdot \infty=\infty \cdot 0=0$. Note that the integral may equal $\infty$.
6.2. Definitions. For an arbitrary function $f: X \rightarrow \overline{\mathbb{R}}$, we define

$$
\begin{aligned}
f^{+}(x) & :=f(x) \text { if } f(x) \geq 0 \\
f^{-}(x) & :=-f(x) \text { if } f(x) \geq 0
\end{aligned}
$$

Thus, $f=f^{+}-f^{-}$and $|f|=f^{+}+f^{-}$.
If $f$ is a measurable countably-simple function and at least one of $\int_{X} f^{+} d \mu$ or $\int_{X} f^{-} d \mu$ is finite, we define

$$
\int_{X} f d \mu:=\int_{X} f^{+} d \mu-\int_{X} f^{-} d \mu .
$$

If $f: X \rightarrow \overline{\mathbb{R}}$ (not necessarily measurable), we define the upper integral of $f$ by

$$
\bar{\int}_{X} f d \mu:=\inf \left\{\int_{X} g d \mu: g \text { is measurable, countably-simple and } g \geq f \mu \text {-a.e. }\right\}
$$

and the lower integral of $f$ by

$$
\int_{X} f d \mu:=\sup \left\{\int_{X} g d \mu: g \text { is measurable, countably-simple and } g \leq f \mu \text {-a.e. }\right\} .
$$

The integral (with respect to the measure $\mu$ ) of a measurable function $f: X \rightarrow \overline{\mathbb{R}}$ is said to exist if

$$
\int_{X} f d \mu=\int_{X} f d \mu,
$$

in which case we write

$$
\int_{X} f d \mu
$$

for the common value. If this value is finite, $f$ is said to be integrable.
6.3. Remark. Observe that our definition requires $f$ to be measurable if it is to be integrable. See Exercise 2, Section 6.2, which shows that measurability is necessary for a function to be integrable provided the measure $\mu$ is complete.
6.4. Remark. If $f$ is a countably-simple function such that $\int_{X} f^{-} d \mu$ is finite then the definitions immediately imply that the integral of $f$ exists and that

$$
\begin{equation*}
\int_{X} f d \mu=\sum_{i=1}^{\infty} a_{i} \mu\left(f^{-1}\left\{a_{i}\right\}\right) \tag{6.1}
\end{equation*}
$$

where the range of $f=\left\{a_{1}, a_{2}, \ldots\right\}$. Clearly, the integral should not depend on the order in which the terms of (6.1) appear. Consequently, the series converges unconditionally, possibly to $+\infty$ (see Exercise 1, Section 6.1). An analogous statement holds if $\int_{X} f^{+} d \mu$ is finite.
6.5. Remark. It is clear from the definitions of upper and lower integrals that if $f=g \mu$-a.e., then $\bar{\int}_{X} f d \mu=\bar{\int}_{X} g d \mu$ and $\underline{\int}_{X} f d \mu=\underline{\int}_{X} g d \mu$. From this observation it follows that if both $f$ and $g$ are measurable, $f=g \mu$-a.e., and $f$ is integrable then $g$ is integrable and

$$
\int_{X} f d \mu=\int_{X} g d \mu
$$

6.6. Definition. If $A \subset X$ (possibly nonmeasurable), we write

$$
\bar{\int}_{A} f d \mu:=\int_{X} f \chi_{A} d \mu
$$

and use analogous notation for the other integrals.

### 6.7. Theorem.

(i) If $f$ is an integrable function, then $f$ is finite $\mu$-a.e.
(ii) If $f$ and $g$ are integrable functions and $a, b$ are constants, then $a f+b g$ is integrable and

$$
\int_{X}(a f+b g) d \mu=a \int_{X} f d \mu+b \int_{X} g d \mu
$$

(iii) If $f$ and $g$ are integrable functions and $f \leq g \mu$-a.e., then

$$
\int_{X} f d \mu \leq \int_{X} g d \mu
$$

(iv) If $f$ is a integrable function and $E \in \mathcal{M}$, then $f \chi_{E}$ is integrable.
(v) A measurable function $f$ is integrable if and only if $|f|$ is integrable.
(vi) If $f$ is a integrable function, then

$$
\left|\int_{X} f d \mu\right| \leq \int_{X}|f| d \mu
$$

Proof. Each of the assertions above is easily seen to hold in case the functions are countably-simple. We leave these proofs as exercises.
(i) If $f$ is integrable, then there are integrable countably-simple functions $g$ and $h$ such that $g \leq f \leq h \mu$-a.e. Thus $f$ is finite $\mu$-a.e.
(ii) Suppose $f$ is integrable and $c$ is a constant. If $c>0$, then for any integrable countably-simple function $g$

$$
c g \leq c f \text { if and only if } g \leq f
$$

Since $\int_{X} c g d \mu=c \int_{X} g d \mu$ it follows that

$$
\underline{\int}_{X} c f d \mu=c \int_{X} f d \mu
$$

and

$$
\bar{\int}_{X} c f d \mu=c \int_{X} f d \mu
$$

Clearly $-f$ is integrable and $\int_{X}-f d \mu=-\int_{X} f d \mu$. Thus if $c<0$, then $c f=|c|(-f)$ is integrable and

$$
\int_{X} c f d \mu=|c| \int_{X}(-f) d \mu=-|c| \int_{X} f d \mu=c \int_{X} f d \mu .
$$

Now suppose $f, g$ are integrable and $f_{1}, g_{1}$ are integrable countably-simple functions such that $f_{1} \leq f, g_{1} \leq g \mu$-a.e. Then $f_{1}+g_{1} \leq f+g \mu$-a.e. and

$$
\underline{\int}_{X}(f+g) d \mu \geq \int_{X}\left(f_{1}+g_{1}\right) d \mu=\int_{X} f_{1} d \mu+\int_{X} g_{1} d \mu .
$$

Thus

$$
\int_{X} f d \mu+\int_{X} g d \mu \leq \int_{X}(f+g) d \mu .
$$

An analogous argument shows

$$
\int_{X}(f+g) d \mu \leq \int_{X} f d \mu+\int_{X} g d \mu
$$

and assertion (ii) follows.
(iii) If $f, g$ are integrable and $f \leq g \mu$-a.e., then, by (ii), $g-f$ is integrable and $g-f \geq 0 \mu$-a.e. . Clearly $\int_{X}(g-f) d \mu=\bar{\int}_{X}(g-f) d \mu \geq 0$ and hence, by (ii) again,

$$
\int_{X} g d \mu=\int_{X} f d \mu+\int_{X}(g-f) d \mu \geq \int_{X} f d \mu .
$$

(iv) If $f$ is integrable, then given $\varepsilon>0$ there are integrable countably-simple functions $g, h$ such that $g \leq f \leq h \mu$-a.e. and

$$
\int_{X}(h-g) d \mu<\varepsilon
$$

Thus

$$
\int_{X}(h-g) \chi_{E} d \mu \leq \varepsilon
$$

for $E \in \mathcal{M}$. Thus

$$
0 \leq \bar{\int}_{X} f \chi_{E} d \mu-\int_{X} f \chi_{E} d \mu<\varepsilon
$$

and since

$$
-\infty<\int_{X} g \chi_{E} d \mu \leq \int_{X} f \chi_{E} d \mu \leq \int_{X} f \chi_{E} d \mu \leq \int_{X} h \chi_{E} d \mu<\infty
$$

it follows that $f \chi_{E}$ is integrable.
(v) If $f$ is integrable, then by (iv) $f^{+}:=f \chi_{\{x: f(x)>0\}}$ and $f^{-}:=-f \chi_{\{x: f(x)<0\}}$ are integrable and by (ii) $|f|=f^{+}+f^{-}$is integrable. If $|f|$ is integrable, then, by (iv), $f^{+}=|f| \chi_{\{x: f(x)>0\}}$ and $f^{-}=|f| \chi_{\{x: f(x)<0\}}$ are integrable and hence $f=f^{+}-f^{-}$is integrable.
(vi) If $f$ is integrable, then by (v) $f^{ \pm}$are integrable and

$$
\left|\int_{X} f d \mu\right|=\left|\int_{X} f^{+} d \mu-\int_{X} f^{-} d \mu\right| \leq \int_{X} f^{+} d \mu+\int_{X} f^{-} d \mu=\int_{X}|f| d \mu
$$

The next result, whose proof is very simple, is remarkably strong in view of the weak hypothesis. In particular, it implies that any bounded, nonnegative, measurable function is $\mu$-integrable. This exhibits a striking difference between the Lebesgue and Riemann integral (see Theorem 6.19 below).
6.8. Theorem. If $f$ is $\mu$-measurable and $f \geq 0 \mu$-a.e., then the integral of $f$ exists: that is,

$$
\int_{X} f d \mu=\bar{\int}_{X} f d \mu
$$

Proof. If the lower integral is infinite, then the upper and lower integrals are both infinite. Thus we may assume that the lower integral is finite and, in particular, $\mu(\{x: f(x)=\infty\})=0$. For $t>1$ and $k=0, \pm 1, \pm 2, \ldots$ set

$$
E_{k}=\left\{x: t^{k} \leq f(x)<t^{k+1}\right\}
$$

and

$$
g_{t}=\sum_{k=-\infty}^{\infty} t^{k} \chi_{E_{k}}
$$

Since each set $E_{k}$ is measurable, it follows that $g_{t}$ is a measurable countably-simple function and $g_{t} \leq f \leq t g_{t} \mu$-a.e. Thus

$$
\bar{\int}_{X} f d \mu \leq \int_{X} t g_{t} d \mu=t \int_{X} g_{t} d \mu \leq t \int_{X} f d \mu
$$

for each $t>1$ and therefore on letting $t \rightarrow 1^{+}$,

$$
\bar{\int}_{X} f d \mu \leq \int_{X} f d \mu
$$

which implies our conclusion since

$$
\int_{X} f d \mu \leq \int_{X} f d \mu
$$

is always true.
6.9. THEOREM. If $f$ is a nonnegative measurable function and $g$ is a integrable function, then

$$
\underline{\int}_{X}(f+g) d \mu=\bar{\int}_{X}(f+g) d \mu=\int_{X} f d \mu+\int_{X} g d \mu .
$$

Proof. If $f$ is integrable, the assertion follows from Theorem 6.7, so assume that $\int_{X} f d \mu=\infty$. Let $h$ be a countably-simple function such that $0 \leq h \leq f$, and let $k$ be an integrable countably-simple function such that $k \geq|g|$. Then, by Exercise 6.1

$$
\underline{\int}_{X}(f+g) d \mu \geq \underline{\int}_{X}(f-|g|) d \mu \geq \int_{X}(h-k) d \mu=\int_{X} h d \mu-\int_{X} k d \mu
$$

from which it follows that $\underline{\int}_{X}(f+g) d \mu=\infty$ and the assertion is proved.
One of the main applications of this result is the following.
6.10. Corollary. If $f$ is measurable, and if either $f^{+}$or $f^{-}$is integrable, then the integral exists:

$$
\int_{X} f d \mu
$$

Proof. For example, if $f^{+}$is integrable, take $g:=-f^{+}$and $f:=f^{-}$in the previous theorem to conclude that the integrals exist:

$$
\int_{X}\left(f^{-}-f^{+}\right) d \mu=\int_{X}-f d \mu
$$

and therefore that

$$
\int_{X} f d \mu
$$

exists.
6.11. ThEOREM. If $f$ is $\mu$-measurable, $g$ is integrable, and $|f| \leq|g| \mu$-a.e. then $f$ is integrable.

Proof. This follows immediately from Theorem 6.7 (v) and Lemma 6.8.

## Exercises for Section 6.1

1. A series $\sum_{i=1}^{\infty} c_{i}$ is said to converge unconditionally if it converges, and for any one-to-one mapping $\sigma$ of $\mathbb{N}$ onto $\mathbb{N}$ the series $\sum_{i=1}^{\infty} c_{\sigma(i)}$ converges to the same limit. Verify the assertion in Remark (6.4). That is, suppose $\mathbb{N}_{1}$ and $\mathbb{N}_{2}$ are both infinite subsets of $\mathbb{N}$ such that $\mathbb{N}_{1} \cap \mathbb{N}_{2}=\emptyset$ and $\mathbb{N}_{1} \cup \mathbb{N}_{2}=\mathbb{N}$. Suppose $\left\{a_{i}: i \in \mathbb{N}\right\}$ are real numbers such that $\left\{a_{i}: i \in \mathbb{N}_{1}\right\}$ are all nonpositive and that $\left\{a_{i}: i \in \mathbb{N}_{2}\right\}$ are all positive numbers. If

$$
-\sum_{i \in \mathbb{N}_{1}} a_{i}<\infty \quad \text { and } \quad \sum_{i \in \mathbb{N}_{2}} a_{i}=\infty
$$

prove that

$$
\sum_{\sigma(i) \in \mathbb{N}} a_{\sigma(i)}=\infty
$$

for any bijection $\sigma: \mathbb{N} \rightarrow \mathbb{N}$. Also, show that

$$
\sum_{\sigma(i) \in \mathbb{N}} a_{\sigma(i)}<\infty \quad \text { and } \quad \sum_{i=1}^{\infty}\left|a_{i}\right|<\infty
$$

if $\sum_{i \in \mathbb{N}_{2}} a_{i}<\infty$. Use the assertion to show that if $f$ is a nonnegative countablysimple function and $g$ is an integrable countably-simple function, then

$$
\int_{X}(f+g) d \mu=\int_{X} f d \mu+\int_{X} g d \mu
$$

2. Verify the assertions of Theorem 6.7 for countably-simple functions.
3. Suppose $f$ is a nonnegative measurable function. Show that

$$
\int_{X} f d \mu=\sup \sum_{k=1}^{N}\left(\inf _{x \in E_{k}} f(x)\right) \mu\left(E_{k}\right)
$$

where the supremum is taken over all finite measurable partitions of $X$, i.e., over all finite collections $\left\{E_{k}\right\}_{k=1}^{N}$ of disjoint measurable subsets of $X$ such that $X=\bigcup_{k=1}^{N} E_{k}$.
4. Suppose $f$ is a nonnegative, integrable function with the property that

$$
\int_{X} f d \mu=0
$$

Show that $f=0 \mu$-a.e.
5. Suppose $f$ is an integrable function with the property that

$$
\int_{E} f d \mu=0
$$

whenever $E$ is an $\mu$-measurable set. Show that $f=0 \mu$-a.e.
6. Show that if $f$ is measurable, $g$ is $\mu$-integrable, and $f \geq g$, then $f^{-}$is $\mu$-integrable and

$$
\underline{\int}_{X} f d \mu=\int_{X} f d \mu=\int_{X} f^{+} d \mu-\int_{X} f^{-} d \mu
$$

7. Suppose $(X, \mathcal{M}, \mu)$ is a measure space and $Y \in \mathcal{M}$. Set

$$
\mu_{Y}(E)=\mu(E \cap Y)
$$

for each $E \in \mathcal{M}$. Show that $\mu_{Y}$ is a measure on $(X, \mathcal{M})$ and that

$$
\int_{X} g d \mu_{Y}=\int_{X} g \chi_{Y} d \mu
$$

for each nonnegative measurable function $g$ on $X$.
8. A function $f:(a, b) \rightarrow \mathbb{R}$ is convex if

$$
f[(1-t) x+t y] \leq(1-t) f(x)+t f(y)
$$

for all $x, y \in(a, b)$ and $t \in[0,1]$. Prove that this is equivalent to

$$
\frac{f(y)-f(x)}{y-x} \leq \frac{f(z)-f(y)}{z-y}
$$

whenever $a<x<y<z<b$.

### 6.2. Limit Theorems

The most important results in integration theory are those related to the continuity of the integral operator. That is, if $\left\{f_{i}\right\}$ converges to $f$ in some sense, how are $\int f$ and $\lim _{i \rightarrow \infty} \int f_{i}$ related? There are three fundamental results that address this question: Fatou's lemma, the Monotone Convergence Theorem, and Lebesgue's Dominated Convergence Theorem. These will be discussed along with associated results.

Our first result concerning the behavior of sequences of integrals is Fatou's Lemma. Note the similarity between this result and its measure-theoretic counterpart, Theorem 4.49, (vi).

We continue to assume the context of a general measure space $(X, \mathcal{M}, \mu)$.
6.12. Lemma (Fatou's Lemma). If $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a sequence of nonnegative $\mu$ measurable functions, then

$$
\int_{X} \liminf _{k \rightarrow \infty} f_{k} d \mu \leq \liminf _{k \rightarrow \infty} \int_{X} f_{k} d \mu .
$$

Proof. Let $g$ be any measurable countably-simple function such that $0 \leq g \leq$ $\liminf _{k \rightarrow \infty} f_{k} \mu$-a.e. . For each $x \in X$, set

$$
g_{k}(x)=\inf \left\{f_{m}(x): m \geq k\right\}
$$

and observe that $g_{k} \leq g_{k+1}$ and

$$
\lim _{k \rightarrow \infty} g_{k}=\liminf _{k \rightarrow \infty} f_{k} \geq g \mu \text {-a.e. }
$$

Write $g=\sum_{j=1}^{\infty} a_{j} \chi_{A_{j}}$ where $A_{j}:=g^{-1}\left(a_{j}\right)$. Therefore $A_{j} \cap A_{i}=\emptyset$ if $i \neq j$ and $X=\bigcup_{k=1}^{\infty} A_{j}$. For $0<t<1$ set

$$
B_{j, k}=A_{j} \cap\left\{x: g_{k}(x)>t a_{j}\right\}
$$

Then each $B_{j, k} \in \mathcal{M}, B_{j, k} \subset B_{j, k+1}$ and since $\lim _{k \rightarrow \infty} g_{k} \geq g \mu$-a.e., we have

$$
\bigcup_{k=1}^{\infty} B_{j, k}=A_{j} \text { and } \lim _{k \rightarrow \infty} \mu\left(B_{j, k}\right)=\mu\left(A_{j}\right)
$$

for $j=1,2, \ldots$ Noting that

$$
\sum_{j=1}^{\infty} t a_{j} \chi_{B_{j, k}} \leq g_{k} \leq f_{m}
$$

for each $m \geq k$, we find

$$
t \int_{X} g d \mu=\sum_{j=1}^{\infty} t a_{j} \mu\left(A_{j}\right)=\lim _{k \rightarrow \infty} \sum_{j=1}^{\infty} t a_{j} \mu\left(B_{j, k}\right) \leq \liminf _{k \rightarrow \infty} \int_{X} f_{k} d \mu
$$

Thus, on letting $t \rightarrow 1^{-}$, we obtain

$$
\int_{X} g d \mu \leq \liminf _{k \rightarrow \infty} \int_{X} f_{k} d \mu
$$

By taking the supremum of the left-hand side over all countable simple functions $g$ with $g \leq \lim \inf _{k \rightarrow \infty} f_{k}$, we have

$$
\int_{X} \liminf _{k \rightarrow \infty} f_{k} d \mu \leq \liminf _{k \rightarrow \infty} \int_{X} f_{k} d \mu
$$

Since $\liminf _{k \rightarrow \infty} f_{k}$ is a nonnegative measurable function, we can apply Theorem 6.8 to obtain our desired conclusion.
6.13. Theorem (Monotone Convergence Theorem). If $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a sequence of nonnegative $\mu$-measurable functions such that $f_{k} \leq f_{k+1}$ for $k=1,2, \ldots$, then

$$
\lim _{k \rightarrow \infty} \int_{X} f_{k} d \mu=\int_{X} \lim _{k \rightarrow \infty} f_{k} d \mu
$$

Proof. Set $f=\lim _{k \rightarrow \infty} f_{k}$. Then $f$ is $\mu$-measurable,

$$
\int_{X} f_{k} d \mu \leq \int_{X} f d \mu \quad \text { for } \quad k=1,2, \ldots
$$

and

$$
\lim _{k \rightarrow \infty} \int_{X} f_{k} d \mu \leq \int_{X} f d \mu
$$

The opposite inequality follows from Fatou's Lemma.
6.14. THEOREM. If $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a sequence of nonnegative $\mu$-measurable functions, then

$$
\int_{X} \sum_{k=1}^{\infty} f_{k} d \mu=\sum_{k=1}^{\infty} \int_{X} f_{k} d \mu
$$

Proof. With $g_{m}:=\sum_{k=1}^{m} f_{k}$ we have $g_{m} \uparrow \sum_{k=1}^{\infty} f_{k}$ and the conclusion follows easily from the Monotone Convergence Theorem and Theorem 6.7 (ii).
6.15. Theorem. If $f$ is integrable and $\left\{E_{k}\right\}_{k=1}^{\infty}$ is a sequence of disjoint measurable sets such that $X=\bigcup_{k=1}^{\infty} E_{k}$, then

$$
\int_{X} f d \mu=\sum_{k=1}^{\infty} \int_{E_{k}} f d \mu
$$

Proof. Assume first that $f \geq 0$, set $f_{k}=f \chi_{E_{k}}$, and apply the previous Theorem. For arbitrary integrable $f$ use the fact that $f=f^{+}-f^{-}$.
6.16. Corollary. If $f \geq 0$ is integrable and if $\nu$ is a set function defined by

$$
\nu(E):=\int_{E} f d \mu
$$

for every measurable set $E$, then $\nu$ is a measure.
6.17. Theorem (Lebesgue's Dominated Convergence Theorem). Suppose $g$ is integrable, $f$ is measurable, $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a sequence of $\mu$-measurable functions such that $\left|f_{k}\right| \leq g \mu$-a.e. for $k=1,2, \ldots$ and

$$
\lim _{k \rightarrow \infty} f_{k}(x)=f(x)
$$

for $\mu$-a.e. $x \in X$. Then

$$
\lim _{k \rightarrow \infty} \int_{X}\left|f_{k}-f\right| d \mu=0
$$

Proof. Clearly $|f| \leq g \mu$-a.e. and hence $f$ and each $f_{k}$ are integrable. Set $h_{k}=2 g-\left|f_{k}-f\right|$. Then $h_{k} \geq 0 \mu$-a.e. and by Fatou's Lemma

$$
\begin{aligned}
2 \int_{X} g d \mu=\int_{X} \liminf _{k \rightarrow \infty} h_{k} d \mu & \leq \liminf _{k \rightarrow \infty} \int_{X} h_{k} d \mu \\
& =2 \int_{X} g d \mu-\limsup _{k \rightarrow \infty} \int_{X}\left|f_{k}-f\right| d \mu
\end{aligned}
$$

Thus

$$
\limsup _{k \rightarrow \infty} \int_{X}\left|f_{k}-f\right| d \mu=0
$$

## Exercises for Section 6.2

1. Let $(X, \mathcal{M}, \mu)$ be an arbitrary measure space. For an arbitrary $X \xrightarrow{f} \overline{\mathbb{R}}$ prove that there is a measurable function with $g \geq f \mu$-a.e such that

$$
\int_{X} g d \mu=\int_{X} f d \mu
$$

2. Suppose $(X, \mathcal{M}, \mu)$ is a measure space, $f: X \rightarrow \overline{\mathbb{R}}$, and

$$
\underline{\int}_{X} f d \mu=\int_{X} f d \mu<\infty
$$

Show that there exists an integrable (and measurable) function $g$ such that $f=g$ $\mu$-a.e. Thus, if $(X, \mathcal{M}, \mu)$ is complete, $f$ is measurable.
3. Suppose $\left\{f_{k}\right\}$ is a sequence of measurable functions, $g$ is a $\mu$-integrable function, and $f_{k} \geq g \mu$-a.e. for each $k$. Show that

$$
\int_{X} \liminf _{k \rightarrow \infty} f_{k} d \mu \leq \liminf _{k \rightarrow \infty} \int_{X} f_{k} d \mu
$$

4. Let $(X, \mathcal{M}, \mu)$ be an arbitrary measure space. For arbitrary nonnegative functions $f_{i}: X \rightarrow \overline{\mathbb{R}}$, prove that

$$
\int_{X} \liminf _{i \rightarrow \infty} f_{i} d \mu \leq \liminf _{i \rightarrow \infty} \int_{X} f_{i} d \mu
$$

Hint: See Exercise 1, Section 6.2.
5. If $\left\{f_{k}\right\}$ is an increasing sequence of measurable functions, $g$ is $\mu$-integrable, and $f_{k} \geq g \mu$-a.e. for each $k$, show that

$$
\lim _{k \rightarrow \infty} \int_{X} f_{k} d \mu=\int_{X} \lim _{k \rightarrow \infty} f_{k} d \mu
$$

6. Show that there exists a sequence of bounded Lebesgue measurable functions mapping $\mathbb{R}$ into $\mathbb{R}$ such that

$$
\liminf _{i \rightarrow \infty} \int_{\mathbb{R}} f_{i} d \lambda<\int_{\mathbb{R}} \liminf _{i \rightarrow \infty} f_{i} d \lambda
$$

7. Let $f$ be a bounded function on the unit square $Q$ in $\mathbb{R}^{2}$. Suppose for each fixed $y$, that $f$ is a measurable function of $x$. For each $(x, y) \in Q$ let the partial derivative $\frac{\partial f}{\partial y}$ exist. Under the assumption that $\frac{\partial f}{\partial y}$ is bounded in $Q$, prove that

$$
\frac{d}{d y} \int_{0}^{1} f(x, y) d \lambda(x)=\int_{0}^{1} \frac{\partial f}{\partial y} d \lambda(x)
$$

### 6.3. Riemann and Lebesgue Integration-A Comparison

The Riemann and Lebesgue integrals are compared, and it is shown that a bounded function is Riemann integrable if and only if it is continuous almost everywhere.

We first recall the definition and some elementary facts concerning Riemann integration. Suppose $[a, b]$ is a closed interval in $\mathbb{R}$. By a partition $\mathcal{P}$ of $[a, b]$ we mean a finite set of points $\left\{x_{i}\right\}_{i=0}^{m}$ such that $a=x_{0}<x_{1}<\cdots<x_{m}=b$. Let

$$
\|\mathcal{P}\|:=\max \left\{x_{i}-x_{i-1}: 1 \leq i \leq m\right\}
$$

For each $i \in\{1,2, \ldots, m\}$ let $x_{i}^{*}$ be an arbitrary point of the interval $\left[x_{i-1}, x_{i}\right]$. A bounded function $f:[a, b] \rightarrow \mathbb{R}$ is Riemann integrable if

$$
\lim _{\|\mathcal{P}\| \rightarrow 0} \sum_{i=1}^{m} f\left(x_{i}^{*}\right)\left(x_{i}-x_{i-1}\right)
$$

exists, in which case the value is the Riemann integral of $f$ over $[a, b]$, which we will denote by

$$
(R) \int_{a}^{b} f(x) d x
$$

Given a partition $\mathcal{P}=\left\{x_{i}\right\}_{i=0}^{m}$ of $[a, b]$ set

$$
\begin{aligned}
U(\mathcal{P}) & =\sum_{i=1}^{m}\left[\sup _{x \in\left[x_{i-1}, x_{i}\right]} f(x)\right]\left(x_{i}-x_{i-1}\right) \\
L(\mathcal{P}) & =\sum_{i=1}^{m}\left[\inf _{x \in\left[x_{i-1}, x_{i}\right]} f(x)\right]\left(x_{i}-x_{i-1}\right) .
\end{aligned}
$$

Then

$$
L(\mathcal{P}) \leq \sum_{i=1}^{m} f\left(x_{i}^{*}\right)\left(x_{i}-x_{i-1}\right) \leq U(\mathcal{P})
$$

for any choice of the $x_{i}^{*}$. Since the supremum (infimum) of $\sum_{i=1}^{m} f\left(x_{i}^{*}\right)\left(x_{i}-x_{i-1}\right)$ over all choices of the $x_{i}^{*}$ is equal to $U(\mathcal{P})(L(\mathcal{P}))$ we see that a bounded function $f$ is Riemann integrable if and only if

$$
\begin{equation*}
\lim _{\|\mathcal{P}\| \rightarrow 0}(U(\mathcal{P})-L(\mathcal{P}))=0 \tag{6.2}
\end{equation*}
$$

We next examine the effects of using a finer partition. Suppose then, that $\mathcal{P}=\left\{x_{i}\right\}_{i=1}^{m}$ is a partition of $[a, b], z \in[a, b]-\mathcal{P}$, and $Q=\mathcal{P} \cup\{z\}$. Thus $\mathcal{P} \subset \mathcal{Q}$ and $\mathcal{Q}$ is called a refinement $\mathcal{P}$. Then $z \in\left(x_{i-1}, x_{i}\right)$ for some $1 \leq i \leq m$ and

$$
\begin{aligned}
\sup _{x \in\left[x_{i-1}, z\right]} f(x) & \leq \sup _{x \in\left[x_{i-1}, x_{i}\right]} f(x) \\
\sup _{x \in\left[z, x_{i}\right]} f(x) & \leq \sup _{x \in\left[x_{i-1}, x_{i}\right]} f(x)
\end{aligned}
$$

Thus

$$
U(\mathcal{Q}) \leq U(\mathcal{P})
$$

An analogous argument shows that $L(\mathcal{P}) \leq L(\mathcal{Q})$. It follows by induction on the number of points in $Q$ that

$$
L(\mathcal{P}) \leq L(\mathcal{Q}) \leq U(\mathcal{Q}) \leq U(\mathcal{P})
$$

whenever $\mathcal{P} \subset \mathcal{Q}$. Thus, $U$ does not increase and $L$ does not decrease when a refinement of the partition is used.

We will say that a Lebesgue measurable function $f$ on $[a, b]$ is Lebesgue integrable if $f$ is integrable with respect to Lebesgue measure $\lambda$ on $[a, b]$.
6.18. Theorem. If $f:[a, b] \rightarrow \mathbb{R}$ is a bounded Riemann integrable function, then $f$ is Lebesgue integrable and

$$
(R) \int_{a}^{b} f(x) d x=\int_{[a, b]} f d \lambda
$$

Proof. Let $\left\{\mathcal{P}_{k}\right\}_{k=1}^{\infty}$ be a sequence of partitions of $[a, b]$ such that $\mathcal{P}_{k} \subset \mathcal{P}_{k+1}$ and $\left\|\mathcal{P}_{k}\right\| \rightarrow 0$ as $k \rightarrow \infty$. Write $\mathcal{P}_{k}=\left\{x_{j}^{k}\right\}_{j=0}^{m_{k}}$. For each $k$ define functions $l_{k}, u_{k}$ by setting

$$
\begin{aligned}
& l_{k}(x)=\inf _{t \in\left[x_{i-1}^{k}, x_{i}^{k}\right]} f(t) \\
& u_{k}(x)=\sup _{t \in\left[x_{i-1}^{k}, x_{i}^{k}\right]} f(t)
\end{aligned}
$$

whenever $x \in\left[x_{i-1}^{k}, x_{i}^{k}\right), 1 \leq i \leq m_{k}$. Then for each $k$ the functions $l_{k}, u_{k}$ are Lebesgue integrable and

$$
\int_{[a, b]} l_{k} d \lambda=L\left(\mathcal{P}_{k}\right) \leq U\left(\mathcal{P}_{k}\right)=\int_{[a, b]} u_{k} d \lambda
$$

The sequence $\left\{l_{k}\right\}$ is monotonically increasing and bounded. Thus

$$
l(x):=\lim _{k \rightarrow \infty} l_{k}(x)
$$

exists for each $x \in[a, b]$ and $l$ is a Lebesgue measurable function. Similarly, the function

$$
u:=\lim _{k \rightarrow \infty} u_{k}
$$

is Lebesgue measurable and $l \leq f \leq u$ on $[a, b]$. Since $f$ is Riemann integrable, it follows from Lebesgue's Dominated Convergence Theorem 6.17 and (6.2) that

$$
\begin{equation*}
\int_{[a, b]}(u-l) d \lambda=\lim _{k \rightarrow \infty} \int_{[a, b]}\left(u_{k}-l_{k}\right) d \lambda=\lim _{k \rightarrow \infty}\left(U\left(\mathcal{P}_{k}\right)-L\left(\mathcal{P}_{k}\right)\right)=0 \tag{6.3}
\end{equation*}
$$

Thus $l=f=u \lambda$-a.e. on $[a, b]$ (see Exercise 4, Section 6.1), and invoking Theorem 6.17 once more we have

$$
\int_{[a, b]} f d \lambda=\lim _{k \rightarrow \infty} \int_{[a, b]} u_{k} d \lambda=\lim _{k \rightarrow \infty} U\left(\mathcal{P}_{k}\right)=(R) \int_{a}^{b} f(x) d x
$$

6.19. Theorem. A bounded function $f:[a, b] \rightarrow \mathbb{R}$ is Riemann integrable if and only if $f$ is continuous $\lambda$-a.e. on $[a, b]$.

Proof. Suppose $f$ is Riemann integrable and let $\left\{\mathcal{P}_{k}\right\}$ be a sequence of partitions of $[a, b]$ such that $\mathcal{P}_{k} \subset \mathcal{P}_{k+1}$ and $\lim _{k \rightarrow \infty}\left\|\mathcal{P}_{k}\right\|=0$. Set $N=\cup_{k=1}^{\infty} \mathcal{P}_{k}$. Let $l_{k}, u_{k}$ be as in the proof of Theorem 6.18. If $x \in[a, b]-N, l(x)=u(x)$ and $\varepsilon>0$, then there is an integer $k$ such that

$$
u_{k}(x)-l_{k}(x)<\varepsilon .
$$

Let $\mathcal{P}_{k}=\left\{x_{j}^{k}\right\}_{j=0}^{m_{k}}$. Then $x \in\left[x_{j-1}^{k}, x_{j}^{k}\right)$ for some $j \in\left\{1,2, \ldots, m_{k}\right\}$. For any $y \in\left(x_{j-1}^{k}, x_{j}^{k}\right)$

$$
|f(y)-f(x)| \leq u_{k}(x)-l_{k}(x)<\varepsilon
$$

Thus $f$ is continuous at $x$. Since $l(x)=u(x)$ for $\lambda$-a.e. $x \in[a, b]$ and $\lambda(N)=0$ we see that $f$ is continuous at $\lambda$-a.e. point of $[a, b]$.

Now suppose $f$ is bounded, $N \subset[a, b]$ with $\lambda(N)=0$, and $f$ is continuous at each point of $[a, b]-N$. Let $\left\{\mathcal{P}_{k}\right\}$ be any sequence of partitions of $[a, b]$ such that $\lim _{k \rightarrow \infty}\left\|\mathcal{P}_{k}\right\|=0$. For each $k$ define the Lebesgue integrable functions $l_{k}$ and $u_{k}$ as in the proof of Theorem 6.18. Then

$$
L\left(\mathcal{P}_{k}\right)=\int_{[a, b]} l_{k} d \lambda \leq \int_{[a, b]} u_{k} d \lambda=U\left(\mathcal{P}_{k}\right)
$$

If $x \in[a, b]-N$ and $\varepsilon>0$, then there is a $\delta>0$ such that

$$
|f(x)-f(y)|<\varepsilon / 2
$$

whenever $|y-x|<\delta$. There is a $k_{0}$ such that $\left\|\mathcal{P}_{k}\right\|<\frac{\delta}{2}$ whenever $k>k_{0}$. Thus

$$
u_{k}(x)-l_{k}(x) \leq 2 \sup \{|f(x)-f(y)|:|y-x|<\delta\}<\varepsilon
$$

whenever $k>k_{0}$. Thus

$$
\lim _{k \rightarrow \infty}\left(u_{k}(x)-l_{k}(x)\right)=0
$$

for each $x \in[a, b]-N$. By the Dominated Convergence Theorem 6.17 it follows that

$$
\lim _{k \rightarrow \infty}\left(U\left(\mathcal{P}_{k}\right)-L\left(\mathcal{P}_{k}\right)\right)=\lim _{k \rightarrow \infty} \int_{[a, b]}\left(u_{k}-l_{k}\right) d \lambda=0
$$

thus showing that $f$ is Riemann integrable.

## Exercises for Section 6.3

1. Give an example of a nondecreasing sequence of functions mapping $[0,1]$ into $[0,1]$ such that each term in the sequence is Riemann integrable and such that the limit of the resulting sequence of Riemann integrals exists, but that the limit of the sequence of functions is not Riemann integrable.
2. From here to Exercise 6 we outline a development of the Riemann-Stieltjes integral that is similar to that of the Riemann integral. Let $f$ and $g$ be two realvalued functions defined on a finite interval $[a, b]$. Given a partition $\mathcal{P}=\left\{x_{i}\right\}_{i=0}^{m}$ of $[a, b]$, for each $i \in\{1,2, \ldots, m\}$ let $x_{i}^{*}$ be an arbitrary point of the interval $\left[x_{i-1}, x_{i}\right]$. We say that the Riemann-Stieltjes integral of $f$ with respect to $g$ exists provided

$$
\lim _{\|\mathcal{P}\| \rightarrow 0} \sum_{i=1}^{m} f\left(x_{i}^{*}\right)\left(g\left(x_{i}\right)-g\left(x_{i-1}\right)\right)
$$

exists, in which case the value is denoted by

$$
\int_{a}^{b} f(x) d g(x)
$$

Prove that if $f$ is continuous and $g$ is continuously differentiable on $[a, b]$, then

$$
\int_{a}^{b} f d g=\int_{a}^{b} f g^{\prime} d x
$$

3. Suppose $f$ is a bounded function on $[\mathrm{a}, \mathrm{b}]$ and $g$ nondecreasing. Set

$$
\begin{aligned}
& U_{R S}(\mathcal{P})=\sum_{i=1}^{m}\left[\sup _{x \in\left[x_{i-1}, x_{i}\right]} f(x)\right]\left(g\left(x_{i}\right)-g\left(x_{i-1}\right)\right) \\
& L_{R S}(\mathcal{P})=\sum_{i=1}^{m}\left[\inf _{x \in\left[x_{i-1}, x_{i}\right]} f(x)\right]\left(g\left(x_{i}\right)-g\left(x_{i-1}\right)\right) .
\end{aligned}
$$

Prove that if $\mathcal{P}^{\prime}$ is a refinement of $\mathcal{P}$, then $L_{R S}\left(\mathcal{P}^{\prime}\right) \geq L_{R S}(\mathcal{P})$ and $U_{R S}\left(\mathcal{P}^{\prime}\right) \leq$ $U_{R S}(\mathcal{P})$. Also, if $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$ are any two partitions, then $L_{R S}\left(\mathcal{P}_{1}\right) \leq U_{R S}\left(\mathcal{P}_{2}\right)$.
4. If $f$ is continuous and $g$ nondecreasing, prove that

$$
\int_{a}^{b} f d g
$$

exists. Thus establish the same conclusion if $g$ is assumed to be of bounded variation.
5. Prove the following integration by parts formula. If $\int_{a}^{b} f d g$ exists, then so does $\int_{a}^{b} g d f$ and

$$
f(b) g(b)-f(a) g(a)=\int_{a}^{b} f d g+\int_{a}^{b} g d f
$$

6. Using the proof of Theorem 6.18 as a guide, show that the Riemann-Stieltjes and Lebesgue-Stieltjes integrals are in agreement. That is, if $f$ is bounded, $g$ is nondecreasing and right-continuous, and if the Riemann-Stieltjes integral of $f$ with respect to $g$ exists, then

$$
\int_{a}^{b} f d g=\int_{[a, b]} f d \lambda_{g}
$$

where $\lambda_{g}$ is the Lebesgue-Stieltjes measure induced by $g$ as in Section 4.6.

### 6.4. Improper Integrals

In this section we study the relation between Lebesgue integrals and improper integrals.

Let $a \in \mathbb{R}$ and $f:[a, \infty) \rightarrow \mathbb{R}$ be a function that is Riemann integrable on each subinterval of $[a, \infty)$. The improper integral of $f$ is defined as

$$
\begin{equation*}
(I) \int_{a}^{\infty} f(x) d x:=\lim _{b \rightarrow \infty}(R) \int_{a}^{b} f(x) d x \tag{6.4}
\end{equation*}
$$

If the limit in (6.4) is finite we say that the improper integral of $f$ exist. We have the following result
6.20. Theorem. Let $f:[a, \infty) \rightarrow \mathbb{R}$ be a nonnegative function that is Riemann integrable on each subinterval of $[a, \infty]$. Then

$$
\begin{equation*}
\int_{[a, \infty)} f d \lambda=\lim _{b \rightarrow \infty}(R) \int_{a}^{b} f d x \tag{6.5}
\end{equation*}
$$

Thus, $f$ is Lebesgue integrable on $[a, \infty)$ if and only if the improper integral (I) $\int_{a}^{\infty} f(x) d x$ exists. Moreover, in this case, $\int_{[a, \infty)} f(x) d \lambda=(I) \int_{a}^{\infty} f(x) d x$.

Proof. Let $b_{n}, n=1,2,3 \ldots$ be any sequence with $b_{n} \rightarrow \infty, b_{n}>a$. We define $f_{n}=f \chi_{\left[a, b_{n}\right]}$. Then the Monotone Convergence Theorem yields

$$
\int_{[a, \infty)} f d \lambda=\lim _{n \rightarrow \infty} \int_{[a, \infty)} f_{n} d \lambda
$$

Thus

$$
\int_{[a, \infty)} f d \lambda=\lim _{n \rightarrow \infty} \int_{\left[a, b_{n}\right]} f d \lambda
$$

Since $f$ is Riemann integrable on each interval $\left[a, b_{n}\right]$, then Theorem 6.18 yields $\int_{\left[a, b_{n}\right]} f d \lambda=(R) \int_{a}^{b_{n}} f d x$ and we conclude

$$
\begin{equation*}
\int_{[a, \infty)} f d \lambda=\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f d x \tag{6.6}
\end{equation*}
$$

The second part follows by noticing that the terms in (6.6) are both finite or $\infty$ at the same time.

If $f:[a, \infty) \rightarrow \mathbb{R}$ takes also negative values then we have the following result:
6.21. Theorem. Let $f:[a, \infty) \rightarrow \mathbb{R}$ be Riemann integrable on every subinterval of $[a, \infty)$. Then $f$ is Lebesgue integrable if and only if the improper integral (I) $\int_{a}^{\infty}|f(x)| d x$ exists. Moreover, in this case,

$$
\begin{equation*}
\int_{[a, \infty)} f d \lambda=(I) \int_{a}^{\infty} f(x) d x \tag{6.7}
\end{equation*}
$$

Proof. Let $f=f^{+}-f^{-}$. Assume that $f$ is Lebesgue integrable on $[a, \infty)$. Thus $f^{+}, f^{-}$are both Lebesgue integrable. From Theorem 6.20 it follows that the improper integrals $(I) \int_{a}^{\infty} f^{+}(x) d x$ and $(I) \int_{a}^{\infty} f^{-}(x) d x$ exist. Moreover, for $b_{n} \rightarrow \infty, b_{n}>a$ we have

$$
\begin{equation*}
\int_{[a, \infty)} f^{+} d \lambda=(I) \int_{a}^{\infty} f^{+} d x=\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{+} d x \tag{6.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{[a, \infty)} f^{-} d \lambda=(I) \int_{a}^{\infty} f^{-} d x=\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{-} d x \tag{6.9}
\end{equation*}
$$

Note that

$$
(R) \int_{a}^{b_{n}} f d x=(R) \int_{a}^{b_{n}} f^{+} d x-(R) \int_{a}^{b_{n}} f^{-} d x
$$

Hence (6.8) and (6.9) imply that

$$
\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f d x=\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{+} d x-\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{-} d x<\infty
$$

which means that the improper integral $(I) \int_{a}^{\infty} f d x$ exists. Moreover (6.8) and (6.9) yield

$$
(I) \int_{a}^{\infty} f d x=\int_{[a, \infty)} f^{+} d \lambda-\int_{[a, \infty)} f^{-} d \lambda=\int_{[a, \infty)} f d \lambda,
$$

which is (6.7). Analogously, using again (6.8) and (6.9),

$$
\begin{equation*}
\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}}|f| d x=\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{+} d x+\lim _{n \rightarrow \infty}(R) \int_{a}^{b_{n}} f^{-} d x<\infty \tag{6.10}
\end{equation*}
$$

and hence the improper integral $(I) \int_{a}^{\infty}|f| d x$ exists with value

$$
(I) \int_{a}^{\infty}|f| d x=\int_{[a, \infty)} f^{+} d \lambda+\int_{[a, \infty)} f^{-} d \lambda=\int_{[a, \infty)}|f| d \lambda .
$$

Conversely, if ( $I$ ) $\int_{a}^{\infty}|f(x)| d x<\infty$ then (6.10) holds and hence by (6.8) and (6.9) we have $\int_{[a, \infty)} f^{+} d \lambda<\infty$ and $\int_{[a, \infty)} f^{-} d \lambda<\infty$, which yield $\int_{[a, \infty)} f d \lambda<\infty$ and hence $f$ is Lebesgue integrable.

## 6.5. $L^{p}$ Spaces

The $L^{p}$ spaces appear in many applications of analysis. They are also the prototypical examples of infinite dimensional Banach spaces which will be studied in Chapter VIII. It will be seen that there is a significant difference in these spaces when $p=1$ and $p>1$.
6.22. Definition. For $1 \leq p \leq \infty$ and $E \in \mathcal{M}$, let $L^{p}(E, \mathcal{M}, \mu)$ denote the class of all measurable functions $f$ on $E$ such that $\|f\|_{p, E ; \mu}<\infty$ where

$$
\|f\|_{p, E ; \mu}:= \begin{cases}\left(\int_{E}|f|^{p} d \mu\right)^{1 / p} & \text { if } 1 \leq p<\infty \\ \inf \{M:|f| \leq M \mu \text {-a.e. on } E\} & \text { if } p=\infty\end{cases}
$$

The quantity $\|f\|_{p, E ; \mu}$ will be called the $L^{p}$ norm of $f$ on $E$ and, for convenience, written $\|f\|_{p}$ when $E=X$ and the measure is clear from the context. The fact that it is a norm will be proved later in this section. We note immediately the following:
(i) $\|f\|_{p} \geq 0$ for any measurable $f$.
(ii) $\|f\|_{p}=0$ if and only if $f=0 \mu$-a.e.
(iii) $\|c f\|_{p}=|c|\|f\|_{p}$ for any $c \in \mathbb{R}$.

For convenience, we will write $L^{p}(X)$ for the class $L^{p}(X, \mathcal{M}, \mu)$. In case $X$ is a topological space, we let $L_{\mathrm{loc}}^{p}(X)$ denote the class of functions $f$ such that $f \in L^{p}(K)$ for each compact set $K \subset X$.

The next lemma shows that the classes $L^{p}(X)$ are vector spaces or, as is more commonly said in this context, linear spaces.
6.23. Theorem. Suppose $1 \leq p \leq \infty$.
(i) If $f, g \in L^{p}(X)$, then $f+g \in L^{p}(X)$.
(ii) If $f \in L^{p}(X)$ and $c \in \mathbb{R}$, then $c f \in L^{p}(X)$.

Proof. Assertion (ii) follows from property (iii) of the $L^{p}$ norm noted above. In case $p$ is finite, assertion (i) follows from the inequality

$$
\begin{equation*}
|a+b|^{p} \leq 2^{p-1}\left(|a|^{p}+|b|^{p}\right) \tag{6.11}
\end{equation*}
$$

which holds for any $a, b \in \mathbb{R}, 1 \leq p<\infty$. For $p \geq 1$, inequality (6.11) follows from the fact that $t \mapsto t^{p}$ is a convex function on $(0, \infty)$ (note that $t \mapsto|t|$ is also convex according the definition given in exercise 6.8) and therefore

$$
\left(\frac{a+b}{2}\right)^{p} \leq \frac{1}{2}\left(a^{p}+b^{p}\right)
$$

In case $p=\infty$, assertion (i) follows from the triangle inequality $|a+b| \leq|a|+|b|$ since if $|f(x)| \leq M \mu$-a.e. and $|g(x)| \leq N \mu$-a.e., then $|f(x)+g(x)| \leq M+N \mu$ a.e.

To deduce further properties of the $L^{p}$ norms we will use the following arithmetic inequality.
6.24. Lemma. For $a, b \geq 0,1<p<\infty$ and $p^{\prime}$ determined by the equation

$$
\frac{1}{p}+\frac{1}{p^{\prime}}=1
$$

we have

$$
a b \leq \frac{a^{p}}{p}+\frac{b^{p^{\prime}}}{p^{\prime}}
$$

Equality holds if and only if $a^{p}=b^{p^{\prime}}$.

Proof. Recall that $\ln (x)$ is an increasing, strictly concave function on $(0, \infty)$, i.e.,

$$
\ln (\lambda x+(1-\lambda) y)>\lambda \ln (x)+(1-\lambda) \ln (y)
$$

for $x, y, \in(0, \infty), x \neq y$ and $\lambda \in[0,1]$.
Set $x=a^{p}, y=b^{p^{\prime}}$, and $\lambda=\frac{1}{p}$ (thus $\left.(1-\lambda)=\frac{1}{p^{\prime}}\right)$ to obtain

$$
\ln \left(\frac{1}{p} a^{p}+\frac{1}{p^{\prime}} b^{p^{\prime}}\right)>\frac{1}{p} \ln \left(a^{p}\right)+\frac{1}{p^{\prime}} \ln \left(b^{p^{\prime}}\right)=\ln (a b)
$$

Clearly equality holds in this inequality if and only if $a^{p}=b^{p^{\prime}}$.
For $p \in[1, \infty]$ the number $p^{\prime}$ defined by $\frac{1}{p}+\frac{1}{p^{\prime}}=1$ is called the Lebesgue conjugate of $p$. We adopt the convention that $p^{\prime}=\infty$ when $p=1$ and $p^{\prime}=1$ when $p=\infty$.
6.25. Theorem. (Hölder's inequality) If $1 \leq p \leq \infty$ and $f, g$ are measurable functions, then

$$
\int_{X}|f g| d \mu=\int_{X}|f||g| d \mu \leq\|f\|_{p}\|g\|_{p^{\prime}}
$$

Equality holds, for $1<p<\infty$, if and only if

$$
\|f\|_{p}^{p}|g|^{p^{\prime}}=\|g\|_{p^{\prime}}^{p^{\prime}}|f|^{p} \quad \mu \text {-a.e.. }
$$

(Recall the convention that $0 \cdot \infty=\infty \cdot 0=0$.)

Proof. In case $p=1$

$$
\int_{X}|f||g| d \mu \leq\|g\|_{\infty} \int_{X}|f| d \mu=\|g\|_{\infty}\|f\|_{1}
$$

and an analogous inequality holds in case $p=\infty$.
In case $1<p<\infty$ the assertion is clear unless $0<\|f\|_{p},\|g\|_{p^{\prime}}<\infty$. In this case, set

$$
\tilde{f}=\frac{f}{\|f\|_{p}} \quad \text { and } \quad \tilde{g}=\frac{g}{\|g\|_{p^{\prime}}}
$$

so that $\|\tilde{f}\|_{p}=1$ and $\|\tilde{g}\|_{p^{\prime}}=1$ and apply Lemma 6.24 to obtain

$$
\frac{1}{\|f\|_{p}\|g\|_{p^{\prime}}} \int_{X}|f||g| d \mu=\int_{X}|\tilde{f}||\tilde{g}| d \mu \leq \frac{1}{p}\|\tilde{f}\|_{p}^{p}+\frac{1}{p^{\prime}}\|\tilde{g}\|_{p^{\prime}}^{p^{\prime}}
$$

The statement concerning equality follows immediately from the preceding lemma when $1<p<\infty$.
6.26. Theorem. Suppose $(X, \mathcal{M}, \mu)$ is a $\sigma$-finite measure space. If $f$ is measurable, $1 \leq p \leq \infty$, and $\frac{1}{p}+\frac{1}{p^{\prime}}=1$, then

$$
\begin{equation*}
\|f\|_{p}=\sup \left\{\int_{X} f g d \mu:\|g\|_{p^{\prime}} \leq 1\right\} \tag{6.12}
\end{equation*}
$$

Proof. Suppose $f$ is measurable. If $g \in L^{p^{\prime}}(X)$ with $\|g\|_{p^{\prime}} \leq 1$, then by Hölder's inequality

$$
\int_{X} f g d \mu \leq\|f\|_{p}\|g\|_{p^{\prime}} \leq\|f\|_{p}
$$

Thus

$$
\sup \left\{\int_{X} f g d \mu:\|g\|_{p^{\prime}} \leq 1\right\} \leq\|f\|_{p}
$$

and it remains to prove the opposite inequality.
In case $p=1$, set $g=\operatorname{sign}(f)$, then $\|g\|_{\infty} \leq 1$ and

$$
\int_{X} f g d \mu=\int_{X}|f| d \mu=\|f\|_{1}
$$

Now consider the case $1<p<\infty$. If $\|f\|_{p}=0$, then $f=0$ a.e. and the desired inequality is clear. If $0<\|f\|_{p}<\infty$, set

$$
g=\frac{|f|^{p / p^{\prime}} \operatorname{sign}(f)}{\|f\|_{p}^{p / p^{\prime}}}
$$

Then $\|g\|_{p^{\prime}}=1$ and

$$
\int f g d \mu=\frac{1}{\|f\|_{p}^{p / p^{\prime}}} \int|f|^{p / p^{\prime}+1} d \mu=\frac{\|f\|_{p}^{p}}{\|f\|_{p}^{p / p^{\prime}}}=\|f\|_{p}
$$

If $\|f\|_{p}=\infty$, let $\left\{X_{k}\right\}_{k=1}^{\infty}$ be an increasing sequence of measurable sets such that $\mu\left(X_{k}\right)<\infty$ for each $k$ and $X=\cup_{k=1}^{\infty} X_{k}$. For each $k$ set

$$
h_{k}(x)=\chi_{X_{k}} \min (|f(x)|, k)
$$

for $x \in X$. Then $h_{k} \in L^{p}(X), h_{k} \leq h_{k+1}$, and $\lim _{k \rightarrow \infty} h_{k}=|f|$. By the Monotone Convergence Theorem $\lim _{k \rightarrow \infty}\left\|h_{k}\right\|_{p}=\infty$. Since we may assume without loss of generality that $\left\|h_{k}\right\|_{p}>0$ for each $k$, there exist, by the result just proved, $g_{k} \in L^{p^{\prime}}(X)$ such that $\left\|g_{k}\right\|_{p^{\prime}}=1$ and

$$
\int h_{k} g_{k} d \mu=\left\|h_{k}\right\|_{p}
$$

Since $h_{k} \geq 0$ we have $g_{k} \geq 0$ and hence

$$
\int f\left(\operatorname{sign}(f) g_{k}\right) d \mu=\int|f| g_{k} d \mu \geq \int h_{k} g_{k} d \mu=\left\|h_{k}\right\|_{p} \rightarrow \infty
$$

as $k \rightarrow \infty$. Thus

$$
\sup \left\{\int f g d \mu:\|g\|_{p^{\prime}} \leq 1\right\}=\infty=\|f\|_{p}
$$

Finally, for the case $p=\infty$, suppose $M:=\sup \left\{\int f g d \mu:\|g\|_{1} \leq 1\right\}<\|f\|_{\infty}$. Thus there exists $\varepsilon>0$ such that $0<M+\varepsilon<\|f\|_{\infty}$. Then the set $E_{\varepsilon}:=\{x:$ $|f(x)| \geq M+\varepsilon\}$ has positive measure, since otherwise we would have $\|f\|_{\infty} \leq M+\varepsilon$. Since $\mu$ is $\sigma$-finite there is a measurable set $E$ such that $0<\mu\left(E_{\varepsilon} \cap E\right)<\infty$. Set

$$
g_{\varepsilon}=\frac{1}{\mu\left(E_{\varepsilon} \cap E\right)} \chi_{E_{\varepsilon} \cap E} \operatorname{sign}(f)
$$

Then $\left\|g_{\varepsilon}\right\|_{1}=1$ and

$$
\int f g_{\varepsilon} d \mu=\frac{1}{\mu\left(E_{\varepsilon} \cap E\right)} \int_{E_{\varepsilon} \cap E}|f| d \mu \geq M+\varepsilon
$$

Thus, $M+\varepsilon \leq \sup \left\{f g d \mu:\|g\|_{1} \leq 1\right\} \leq\|f\|_{\infty}$ which contradicts that $\sup \{f g d \mu$ : $\left.\|g\|_{1} \leq 1\right\}=M$. We conclude

$$
\sup \left\{\int f g d \mu:\|g\|_{1} \leq 1\right\}=\|f\|_{\infty}
$$

6.27. Theorem (Minkowski's inequality). Suppose $1 \leq p \leq \infty$ and $f, g \in$ $L^{p}(X)$. Then

$$
\|f+g\|_{p} \leq\|f\|_{p}+\|g\|_{p}
$$

Proof. The assertion is clear in case $p=1$ or $p=\infty$, so suppose $1<p<\infty$. Then, applying first the triangle inequality and then Hölder's inequality, we obtain

$$
\begin{aligned}
\|f+g\|_{p}^{p}= & \int|f+g|^{p} d \mu=\int|f+g|^{p-1}|f+g| \\
\leq & \int|f+g|^{p-1}|f| d \mu+\int|f+g|^{p-1}|g| d \mu \\
\leq & \left(\int\left(|f+g|^{p-1}\right)^{p^{\prime}}\right)^{1 / p^{\prime}}\left(\int|f|^{p} d \mu\right)^{1 / p} \\
& +\left(\int\left(|f+g|^{p-1}\right)^{p^{\prime}}\right)^{1 / p^{\prime}}\left(\int|g|^{p} d \mu\right)^{1 / p} \\
\leq & \left(\int|f+g|^{p}\right)^{(p-1) / p}\left(\int|f|^{p} d \mu\right)^{1 / p} \\
& +\left(\int|f+g|^{p}\right)^{(p-1) / p}\left(\int|g|^{p} d \mu\right)^{1 / p} \\
= & \|f+g\|_{p}^{p-1}\|f\|_{p}+\|f+g\|_{p}^{p-1}\|g\|_{p} \\
= & \|f+g\|_{p}^{p-1}\left(\|f\|_{p}+\|g\|_{p}\right)
\end{aligned}
$$

The assertion is clear if $\|f+g\|_{p}=0$. Otherwise we divide by $\|f+g\|_{p}^{p-1}$ to obtain

$$
\|f+g\|_{p} \leq\|f\|_{p}+\|g\|_{p}
$$

As a consequence of Theorem 6.27 and the remarks following Definition 6.22 we can say that for $1 \leq p \leq \infty$ the spaces $L^{p}(X)$ are, in the terminology of Chapter 8 , normed linear spaces provided we agree to identify functions that are equal $\mu$-a.e. The norm $\|\cdot\|_{p}$ induces a metric $\rho$ on $L^{p}(X)$ if we define

$$
\rho(f, g):=\|f-g\|_{p}
$$

for $f, g \in L^{p}(X)$ and agree to interpret the statement " $f=g$ " as $f=g \mu$-a.e.
6.28. Definitions. A sequence $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a Cauchy sequence in $L^{p}(X)$ if given any $\varepsilon>0$ there is a positive integer $N$ such that

$$
\left\|f_{k}-f_{m}\right\|_{p}<\varepsilon
$$

whenever $k, m>N$. The sequence $\left\{f_{k}\right\}_{k=1}^{\infty}$ converges in $L^{p}(X)$ to $f \in L^{p}(X)$ if

$$
\lim _{k \rightarrow \infty}\left\|f_{k}-f\right\|_{p}=0
$$

6.29. Theorem. If $1 \leq p \leq \infty$, then $L^{p}(X)$ is a complete metric space under the metric $\rho$, i.e., if $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a Cauchy sequence in $L^{p}(X)$, then there is an $f \in L^{p}(X)$ such that

$$
\lim _{k \rightarrow \infty}\left\|f_{k}-f\right\|_{p}=0
$$

Proof. Suppose $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a Cauchy sequence in $L^{p}(X)$. There an integer $N$ such that $\left\|f_{k}-f_{m}\right\|_{p}<1$ whenever $k, m \geq N$. By Minkowski's inequality

$$
\left\|f_{k}\right\|_{p} \leq\left\|f_{N}\right\|_{p}+\left\|f_{k}-f_{N}\right\|_{p} \leq\left\|f_{N}\right\|_{p}+1
$$

whenever $k \geq N$. Thus the sequence $\left\{\left\|f_{k}\right\|_{p}\right\}_{k=1}^{\infty}$ is bounded.
Consider the case $1 \leq p<\infty$. For any $\varepsilon>0$, let $A_{k, m}:=\left\{x:\left|f_{k}(x)-f_{m}(x)\right| \geq\right.$ $\varepsilon\}$. Then,

$$
\int_{A_{k, m}}\left|f_{k}-f_{m}\right|^{p} d \mu \geq \varepsilon^{p} \mu\left(A_{k, m}\right)
$$

that is,

$$
\varepsilon^{p} \mu\left(\left\{x:\left|f_{k}(x)-f_{m}(x)\right| \geq \varepsilon\right\}\right) \leq\left\|f_{k}-f_{m}\right\|_{p}^{p}
$$

Thus $\left\{f_{k}\right\}$ is fundamental in measure, and consequently by Exercise 6, Section 5.2, and Theorem 5.26, there exists a subsequence $\left\{f_{k_{j}}\right\}_{j=1}^{\infty}$ that converges $\mu$-a.e. to a measurable function $f$. By Fatou's lemma,

$$
\|f\|_{p}^{p}=\int|f|^{p} d \mu \leq \liminf _{j \rightarrow \infty} \int\left|f_{k_{j}}\right|^{p} d \mu<\infty
$$

Thus $f \in L^{p}(X)$.
Let $\varepsilon>0$ and let $M$ be such that $\left\|f_{k}-f_{m}\right\|_{p}<\varepsilon$ whenever $k, m>M$. Using Fatou's lemma again we see

$$
\left\|f_{k}-f\right\|_{p}^{p}=\int\left|f_{k}-f\right|^{p} d \mu \leq \liminf _{j \rightarrow \infty} \int\left|f_{k}-f_{k_{j}}\right|^{p} d \mu<\varepsilon^{p}
$$

whenever $k>M$. Thus $f_{k}$ converges to $f$ in $L^{p}(X)$.
The case $p=\infty$ is left as Exercise 2, Section 6.5.
As a consequence of Theorem 6.29 we see for $1 \leq p \leq \infty$, that $L^{p}(X)$ is a Banach space; i.e., a normed linear space that is complete with respect to the metric induced by the norm.

Here we include a useful result relating norm convergence in $L^{p}$ and pointwise convergence.
6.30. Theorem (Vitali's Convergence Theorem). Suppose $\left\{f_{k}\right\}, f \in L^{p}(X)$, $1 \leq p<\infty$. Then $\left\|f_{k}-f\right\|_{p} \rightarrow 0$ if the following three conditions hold:
(i) $f_{k} \rightarrow f \mu$-a.e.
(ii) For each $\varepsilon>0$, there exists a measurable set $E$ such that $\mu(E)<\infty$ and

$$
\int_{\widetilde{E}}\left|f_{k}\right|^{p} d \mu<\varepsilon, \quad \text { for all } \quad k \in \mathbb{N}
$$

(iii) For each $\varepsilon>0$, there exists $\delta>0$ such that $\mu(E)<\delta$ implies

$$
\int_{E}\left|f_{k}\right|^{p} d \mu<\varepsilon \quad \text { for all } \quad k \in \mathbb{N} .
$$

Conversely, if $\left\|f_{k}-f\right\|_{p} \rightarrow 0$, then (ii) and (iii) hold. Furthermore, (i) holds for a subsequence.

Proof. Assume the three conditions hold. Choose $\varepsilon>0$ and let $\delta>0$ be the corresponding number given by (iii). Condition (ii) provides a measurable set $E$ with $\mu(E)<\infty$ such that

$$
\int_{\tilde{E}}\left|f_{k}\right|^{p} d \mu<\varepsilon
$$

for all positive integers $k$. Since $\mu(E)<\infty$, we can apply Egoroff's Theorem to obtain a measurable set $B \subset E$ with $\mu(E-B)<\delta$ such that $f_{k}$ converges uniformly to $f$ on $B$. Now write

$$
\begin{aligned}
\int_{X}\left|f_{k}-f\right|^{p} d \mu & =\int_{B}\left|f_{k}-f\right|^{p} d \mu \\
& +\int_{E-B}\left|f_{k}-f\right|^{p} d \mu+\int_{\tilde{E}}\left|f_{k}-f\right|^{p} d \mu
\end{aligned}
$$

The first integral on the right can be made arbitrarily small for large $k$, because of the uniform convergence of $f_{k}$ to $f$ on $B$. The second and third integrals will be estimated with the help of the inequality

$$
\left|f_{k}-f\right|^{p} \leq 2^{p-1}\left(\left|f_{k}\right|^{p}+|f|^{p}\right)
$$

see (6.11). From (iii) we have $\int_{E-B}\left|f_{k}\right|^{p}<\varepsilon$ for all $k \in \mathbb{N}$ and then Fatou's Lemma shows that $\int_{E-B}|f|^{p}<\varepsilon$ as well. The third integral can be handled in a similar way using (ii). Thus, it follows that $\left\|f_{k}-f\right\|_{p} \rightarrow 0$.

Now suppose $\left\|f_{k}-f\right\|_{p} \rightarrow 0$. Then for each $\varepsilon>0$ there exists a positive integer $k_{0}$ such that $\left\|f_{k}-f\right\|_{p}<\varepsilon / 2$ for $k>k_{0}$. With the help of Exercise 3, Section 6.5, there exist measurable sets $A$ and $B$ of finite measure such that

$$
\int_{\tilde{A}}|f|^{p} d \mu<(\varepsilon / 2)^{p} \quad \text { and } \quad \int_{\tilde{B}}\left|f_{k}\right|^{p} d \mu<(\varepsilon)^{p} \quad \text { for } \quad k=1,2, \ldots, k_{0}
$$

Minkowski's inequality implies that

$$
\left\|f_{k}\right\|_{p, \tilde{A}} \leq\left\|f_{k}-f\right\|_{p, \tilde{A}}+\|f\|_{p, \tilde{A}}<\varepsilon \quad \text { for } \quad k>k_{0}
$$

Then set $E=A \cup B$ to obtain the necessity of (ii).
Similar reasoning establishes the necessity of (iii).

According to Exercise 4, Section 6.5 convergence in $L^{p}$ implies convergence in measure. Hence, (i) holds for a subsequence.

Finally, we conclude this section by considering how $L^{p}(X)$ compares with $L^{q}(X)$ for $1 \leq p<q \leq \infty$. For example, let $X=[0,1]$ and let $\mu:=\lambda$. In this case it is easy to see that $L^{q} \subset L^{p}$, for if $f \in L^{q}$ then $|f(x)|^{q} \geq|f(x)|^{p}$ if $x \in A:=\{x:|f(x)| \geq 1\}$. Therefore,

$$
\int_{A}|f|^{p} d \lambda \leq \int_{A}|f|^{q} d \lambda<\infty
$$

while

$$
\int_{[0,1] \backslash A}|f|^{p} d \lambda \leq 1 \cdot \lambda([0,1])<\infty .
$$

This observation extends to a more general situation via Hölder's inequality.
6.31. Theorem. If $\mu(X)<\infty$ and $1 \leq p \leq q \leq \infty$, then $L^{q}(X) \subset L^{p}(X)$ and

$$
\|f\|_{q ; \mu} \leq \mu(X)^{\frac{1}{p}-\frac{1}{q}}\|f\|_{p ; \mu}
$$

Proof. If $q=\infty$, then result is immediate:

$$
\|f\|_{p}^{p}=\int_{X}|f|^{p} d \mu \leq\|f\|_{\infty}^{p} \int_{X} 1 d \mu=\|f\|_{\infty}^{p} \mu(X)<\infty
$$

If $q<\infty$ Hölder's inequality with conjugate exponents $q / p$ and $q /(q-p)$ implies that

$$
\|f\|_{p}^{p}=\int_{X}|f|^{p} \cdot 1 d \mu \leq\left\||f|^{p}\right\|_{q / p}\|1\|_{q /(q-p)}=\|f\|_{q}^{p} \mu(X)^{(q-p) / q}<\infty .
$$

6.32. Theorem. If $0<p<q<r \leq \infty$, then $L^{p}(X) \cap L^{r}(X) \subset L^{q}(X)$ and

$$
\|f\|_{q ; \mu} \leq\|f\|_{p ; \mu}^{\lambda}\|f\|_{r ; \mu}^{1-\lambda}
$$

where $0<\lambda<1$ is defined by he equation

$$
\frac{1}{q}=\frac{\lambda}{p}+\frac{1-\lambda}{r}
$$

Proof. If $r<\infty$, use Hölder's inequality with conjugate indices $p / \lambda q$ and $r /(1-\lambda) q$ to obtain

$$
\begin{aligned}
\int_{X}|f|^{q}=\int_{X}|f|^{\lambda q}|f|^{(1-\lambda) q} & \leq\left\||f|^{\lambda q}\right\|_{p / \lambda q}\left\||f|^{(1-\lambda) q}\right\|_{r /(1-\lambda) q} \\
& =\left(\int_{X}|f|^{p}\right)^{\lambda q / p}\left(\int_{X}|f|^{r}\right)^{(1-\lambda) q / r} \\
& =\|f\|_{p}^{\lambda q}\|f\|_{r}^{(1-\lambda) q}
\end{aligned}
$$

We obtain the desired result by taking $q^{\text {th }}$ roots of both sides.

When $r=\infty$, we have

$$
\int_{X}|f|^{q} \leq\|f\|_{\infty}^{q-p} \int_{X}|f|^{p}
$$

and so

$$
\|f\|_{q} \leq\|f\|_{p}^{p / q}\|f\|_{\infty}^{1-(p / q)}=\|f\|_{p}^{\lambda}\|f\|_{\infty}^{1-\lambda}
$$

## Exercises for Section 6.5

1. Use Theorem 5.27 to show that if $f \in L^{p}(X)(1 \leq p<\infty)$, then there is a sequence $\left\{f_{k}\right\}$ of measurable simple functions such that $\left|f_{k}\right| \leq|f|$ for each $k$ and

$$
\lim _{k \rightarrow \infty}\left\|f-f_{k}\right\|_{L^{p}(X)}=0
$$

2. Prove Theorem 6.29 in case $p=\infty$.
3. Suppose $(X, \mathcal{M}, \mu)$ is an arbitrary measure space, $\|f\|_{p}<\infty, 1 \leq p<\infty$, and $\varepsilon>0$. Prove that there is a measurable set $E$ with $\mu(E)<\infty$ such that

$$
\int_{\tilde{E}}|f|^{p} d \mu<\varepsilon
$$

4. Prove that convergence in $L^{p}, 1 \leq p<\infty$, implies convergence in measure.
5. Let $(X, \mathcal{M}, \mu)$ be a $\sigma$-finite measure space. Prove that there is a function $f \in$ $L^{1}(\mu)$ such that $0<f<1$ everywhere on $X$.
6. Suppose $\mu$ and $\nu$ are measures on $(X, \mathcal{M})$ with the property that $\mu(E) \leq \nu(E)$ for each $E \in \mathcal{M}$. For $p \geq 1$ and $f \in L^{p}(X, \nu)$, show that $f \in L^{p}(X, \mu)$ and that

$$
\int_{X}|f|^{p} d \mu \leq \int_{X}|f|^{p} d \nu
$$

7. Suppose $f \in L^{p}(X, \mathcal{M}, \mu), 1 \leq p<\infty$. Then for any $t>0$,

$$
\mu(\{|f|>t\}) \leq t^{-p}\|f\|_{p ; \mu}^{p} .
$$

This is known as Chebyshev's Inequality.
8. Prove that a differentiable function $f$ on $(a, b)$ is convex if and only if $f^{\prime}$ is monotonically increasing.
9. Prove that a convex function is continuous.
10. (a) Prove Jensen's inequality: Let $f \in L^{1}(X, \mathcal{M}, \mu)$ where $\mu(X)<\infty$ and suppose $f(X) \subset[a, b]$. If $\varphi$ is a convex function on $[a, b]$, then

$$
\varphi\left(\frac{1}{\mu(X)} \int_{X} f d \mu\right) \leq \frac{1}{\mu(X)} \int_{X}(\varphi \circ f) d \mu
$$

Thus, $\varphi(\operatorname{average}(f)) \leq \operatorname{average}(\varphi \circ f)$. Hint: let

$$
t_{0}=\left[\mu(X)^{-1}\right] \int_{X} f d \mu . \quad \text { Then } \quad t_{0} \in(a, b) .
$$

Furthermore, with

$$
\alpha:=\sup _{t \in\left(a, t_{0}\right)} \frac{\varphi\left(t_{0}\right)-\varphi(t)}{t_{0}-t},
$$

we have $\varphi(t)-\varphi\left(t_{0}\right) \geq \alpha\left(t-t_{0}\right)$ for all $t \in(a, b)$. In particular, $\varphi(f(x))-$ $\varphi\left(t_{0}\right) \geq \alpha\left(f(x)-t_{0}\right)$ for all $x \in X$. Now integrate.
(b) Observe that if $\varphi(t)=t^{p}, 1 \leq p<\infty$, then Jensen's inequality follows from Hölder's inequality:

$$
\begin{gathered}
\Longrightarrow \quad[\mu(X)]^{-1} \int_{X} f \cdot 1 d \mu \leq\|f\|_{p}[\mu(X)]^{\frac{1}{p^{\prime}}-1}=\|f\|_{p}[\mu(X)]^{-1 / p} \\
\\
\quad\left(\frac{1}{\mu(X)} \int_{X} f d \mu\right)^{p} \leq \frac{1}{\mu(X)} \int_{X}\left(|f|^{p}\right) d \mu
\end{gathered}
$$

(c) However, Jensen's inequality is stronger than Hölder's inequality in the following sense: If $f$ is defined on $[0,1]$ then

$$
e^{\int_{X} f d \lambda} \leq \int_{X} e^{f(x)} d \lambda
$$

(d) Suppose $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ is such that

$$
\varphi\left(\int_{0}^{1} f d \lambda\right) \leq \int_{0}^{1} \varphi(f) d \lambda
$$

for every real bounded measurable $f$. Prove that $\varphi$ is convex.
(e) Thus, we have

$$
\varphi\left(\int_{0}^{1} f d \lambda\right) \leq \int_{0}^{1} \varphi(f) d \lambda
$$

for each bounded, measurable $f$ if and only if $\varphi$ is convex.
11. In the context of a measure space $(X, \mathcal{M}, \mu)$, suppose $f$ is a bounded measurable function with $a \leq f(x) \leq b$ for $\mu$-a.e. $x \in X$. Prove that for each integrable function $g$, there exists a number $c \in[a, b]$ such that

$$
\int_{X} f|g| d \mu=c \int_{X}|g| d \mu
$$

12. (a) Suppose $f$ is a Lebesgue integrable function on $\mathbb{R}^{n}$. Prove that for each $\varepsilon>0$ there is a continuous function $g$ with compact support on $\mathbb{R}^{n}$ such that

$$
\int_{\mathbb{R}^{n}}|f(y)-g(y)| d \lambda(y)<\varepsilon
$$

(b) Show that the above result is true for $f \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<\infty$. That is, show that the continuous functions with compact support are dense in $L^{p}\left(\mathbb{R}^{n}\right)$. Hint: Use Corollary 5.28 to show that step functions are dense in $L^{p}\left(\mathbb{R}^{n}\right)$.
13. If $f \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<\infty$, then prove

$$
\lim _{|h| \rightarrow 0}\|f(x+h)-f(x)\|_{p}=0
$$

Also, show that this result fails when $p=\infty$.
14. Let $p_{1}, p_{2}, \ldots, p_{m}$ be positive real numbers such that

$$
\sum_{i=1}^{m} p_{i}=1
$$

For $f_{1}, f_{2}, \ldots, f_{m} \in L^{1}(X, \mu)$, prove that

$$
f_{1}^{p_{1}} f_{2}^{p_{2}} \cdots f_{m}^{p_{m}} \in L^{1}(X, \mu)
$$

and

$$
\int_{X}\left(f_{1}^{p_{1}} f_{2}^{p_{2}} \cdots f_{m}^{p_{m}}\right) d \mu \leq\left\|f_{1}\right\|_{1}^{p_{1}}\left\|f_{2}\right\|_{1}^{p_{2}} \cdots\left\|f_{m}\right\|_{1}^{p_{m}}
$$

### 6.6. Signed Measures

We develop the basic properties of countably additive set functions of arbitrary sign, or signed measures. In particular, we establish the decomposition theorems of Hahn and Jordan, which show that signed measures and (positive) measures are closely related.

Let $(X, \mathcal{M}, \mu)$ be a measure space. Suppose $f$ is measurable, at least one of $f^{+}$, $f^{-}$is integrable, and set

$$
\begin{equation*}
\nu(E)=\int_{E} f d \mu \tag{6.13}
\end{equation*}
$$

for $E \in \mathcal{M}$. (Recall from Corollary 6.10, that the integral in (6.13) exists.) Then $\nu$ is an extended real-valued function on $\mathcal{M}$ with the following properties:
(i) $\nu$ assumes at most one of the values $+\infty,-\infty$,
(ii) $\nu(\emptyset)=0$,
(iii) If $\left\{E_{k}\right\}_{k=1}^{\infty}$ is a disjoint sequence of measurable sets then

$$
\nu\left(\bigcup_{k=1}^{\infty} E_{k}\right)=\sum_{k=1}^{\infty} \nu\left(E_{k}\right)
$$

where the series on the right either converges absolutely or diverges to $\pm \infty$ (see Exercise 1, Section 6.6),
(iv) $\nu(E)=0$ whenever $\mu(E)=0$.

In Section 6.6 we will show that the properties (i)-(iv) characterize set functions of the type (6.13).
6.33. Remark. An extended real-valued function $\nu$ defined on $\mathcal{M}$ is a signed measure if it satisfies properties (i)-(iii) above. If in addition it satisfies (iv) the signed measure $\nu$ is said to be absolutely continuous with respect to $\mu$, written $\nu \ll \mu$. In some contexts, we will underscore that a measure $\mu$ is not a signed measure by saying that it is a positive measure. In other words, a positive measure is merely a measure in the sense defined in Definition 4.47.
6.34. Definition. Let $\nu$ be a signed measure on $\mathcal{M}$. A set $A \in \mathcal{M}$ is a positive set for $\nu$ if $\nu(E) \geq 0$ for each measurable subset $E$ of $A$. A set $B \in \mathcal{M}$ is a negative set for $\nu$ if $\nu(E) \leq 0$ for each measurable subset $E$ of $B$. A set $C \in \mathcal{M}$ is a null set for $\nu$ if $\nu(E)=0$ for each measurable subset $E$ of $C$.

Note that any measurable subset of a positive set for $\nu$ is also a positive set and that analogous statements hold for negative sets and null sets. It follows that any countable union of positive sets is a positive set. To see this suppose $\left\{P_{k}\right\}_{k=1}^{\infty}$ is a sequence of positive sets. Then there exist disjoint measurable sets $P_{k}^{*} \subset P_{k}$ such that $P:=\bigcup_{k=1}^{\infty} P_{k}=\bigcup_{k=1}^{\infty} P_{k}^{*}$ (Lemma 4.7). If $E$ is a measurable subset of $P$, then

$$
\nu(E)=\sum_{k=1}^{\infty} \nu\left(E \cap P_{k}^{*}\right) \geq 0
$$

since each $P_{k}^{*}$ is positive for $\nu$.
It is important to observe the distinction between measurable sets $E$ such that $\nu(E)=0$ and null sets for $\nu$. If $E$ is a null set for $\nu$, then $\nu(E)=0$ but the converse is not generally true.
6.35. Theorem. If $\nu$ is a signed measure on $\mathcal{M}, E \in \mathcal{M}$ and $0<\nu(E)<\infty$ then $E$ contains a positive set $A$ with $\nu(A)>0$.

Proof. If $E$ is positive then the conclusion holds for $A=E$. Assume $E$ is not positive, and inductively construct a sequence of sets $E_{k}$ as follows. Set

$$
c_{1}:=\inf \{\nu(B): B \in \mathcal{M}, B \subset E\}<0
$$

There exists a measurable set $E_{1} \subset E$ such that

$$
\nu\left(E_{1}\right)<\frac{1}{2} \max \left(c_{1},-1\right)<0
$$

For $k \geq 1$, if $E \backslash \bigcup_{j=1}^{k} E_{j}$ is not positive, then

$$
c_{k+1}:=\inf \left\{\nu(B): B \in \mathcal{M}, B \subset E \backslash \bigcup_{j=1}^{k} E_{j}\right\}<0
$$

and there is a measurable set $E_{k+1} \subset E \backslash \bigcup_{j=1}^{k} E_{j}$ such that

$$
\nu\left(E_{k+1}\right)<\frac{1}{2} \max \left(c_{k+1},-1\right)<0
$$

Note that if $c_{k}=-\infty$, then $\nu\left(E_{k}\right)<-1 / 2$.
If at any stage $E \backslash \bigcup_{j=1}^{k} E_{j}$ is a positive set, let $A=E \backslash \bigcup_{j=1}^{k} E_{j}$ and observe

$$
\nu(A)=\nu(E)-\sum_{j=1}^{k} \nu\left(E_{j}\right)>\nu(E)>0
$$

Otherwise set $A=E \backslash \bigcup_{k=1}^{\infty} E_{k}$ and observe

$$
\nu(E)=\nu(A)+\sum_{k=1}^{\infty} \nu\left(E_{k}\right)
$$

Since $\nu(E)>0$ we have $\nu(A)>0$. Since $\nu(E)$ is finite the series converges absolutely, $\nu\left(E_{k}\right) \rightarrow 0$ and therefore $c_{k} \rightarrow 0$ as $k \rightarrow \infty$. If $B$ is a measurable subset of $A$, then $B \bigcap E_{k}=\emptyset$ for $k=1,2, \ldots$ and hence

$$
\nu(B) \geq c_{k}
$$

for $k=1,2, \ldots$. Thus $\nu(B) \geq 0$. This shows that $A$ is a positive set and the lemma is proved.
6.36. Theorem (Hahn Decomposition). If $\nu$ is a signed measure on $\mathcal{M}$, then there exist disjoint sets $P$ and $N$ such that $P$ is a positive set, $N$ is a negative set, and $X=P \cup N$.

Proof. By considering $-\nu$ in place of $\nu$ if necessary, we may assume $\nu(E)<\infty$ for each $E \in \mathcal{M}$. Set

$$
\lambda:=\sup \{\nu(A): A \text { is a positive set for } \nu\} .
$$

Since $\emptyset$ is a positive set, $\lambda \geq 0$. Let $\left\{A_{k}\right\}_{k=1}^{\infty}$ be a sequence of positive sets for which

$$
\lim _{k \rightarrow \infty} \nu\left(A_{k}\right)=\lambda
$$

Set $P=\bigcup_{k=1}^{\infty} A_{k}$. Then $P$ is positive and hence $\nu(P) \leq \lambda$. On the other hand each $P \backslash A_{k}$ is positive and hence

$$
\nu(P)=\nu\left(A_{k}\right)+\nu\left(P \backslash A_{k}\right) \geq \nu\left(A_{k}\right)
$$

Thus $\nu(P)=\lambda<\infty$.

Set $N=X \backslash P$. We have only to show that $N$ is negative. Suppose $B$ is a measurable subset of $N$. If $\nu(B)>0$, then by Lemma $6.35 B$ must contain a positive set $B^{*}$ such that $\nu\left(B^{*}\right)>0$. But then $B^{*} \bigcup P$ is positive and

$$
\nu\left(B^{*} \cup P\right)=\nu\left(B^{*}\right)+\nu(P)>\lambda
$$

contradicting the choice of $\lambda$.
Note that the Hahn decomposition above is not unique if $\nu$ has a nonempty null set.

The following definition describes a relation between measures that is the antithesis of absolute continuity.
6.37. Definition. Two measures $\mu_{1}$ and $\mu_{2}$ defined on a measure space $(X, \mathcal{M})$ are said to be mutually singular (written $\mu_{1} \perp \mu_{2}$ ) if there exists a measurable set $E$ such that

$$
\mu_{1}(E)=0=\mu_{2}(X-E)
$$

6.38. Theorem (Jordan Decomposition). If $\nu$ is a signed measure on $\mathcal{M}$ then there exists a unique pair of mutually singular measures $\nu^{+}$and $\nu^{-}$, at least one of which is finite, such that

$$
\nu(E)=\nu^{+}(E)-\nu^{-}(E)
$$

for each $E \in \mathcal{M}$.
Proof. Let $P \cup N$ be a Hahn decomposition of $X$ with $P \cap N=\emptyset, P$ positive and $N$ negative for $\nu$. Set

$$
\begin{aligned}
& \nu^{+}(E)=\nu(E \cap P) \\
& \nu^{-}(E)=-\nu(E \cap N)
\end{aligned}
$$

for $E \in \mathcal{M}$. Clearly $\nu^{+}$and $\nu^{-}$are measures on $\mathcal{M}$ and $\nu=\nu^{+}-\nu^{-}$. The measures $\nu^{+}$and $\nu^{-}$are mutually singular since $\nu^{+}(N)=0=\nu^{-}(X-N)$. That at least one of the measures $\nu^{+}, \nu^{-}$is finite follows immediately from the fact that $\nu^{+}(X)=\nu(P)$ and $\nu^{-}(X)=-\nu(N)$, at least one of which is finite.

If $\nu_{1}$ and $\nu_{2}$ are positive measures such that $\nu=\nu_{1}-\nu_{2}$ and $A \in \mathcal{M}$ such that $\nu_{1}(X-A)=0=\nu_{2}(A)$, then

$$
\begin{aligned}
\nu_{1}(X-P) & =\nu_{1}((X \cap A) \backslash P) \\
& =\nu((X \cap A) \backslash P)+\nu_{2}((X \cap A) \backslash P) \\
& =-\nu^{-}(X \cap A) \leq 0
\end{aligned}
$$

Thus $\nu_{1}(X \backslash P)=0$. Similarly $\nu_{2}(P)=0$. For any $E \in \mathcal{M}$ we have

$$
\nu^{+}(E)=\nu(E \cap P)=\nu_{1}(E \cap P)-\nu_{2}(E \cap P)=\nu_{1}(E) .
$$

Analogously $\nu^{-}=\nu_{2}$.
Note that if $\nu$ is the signed measure defined by (6.13) then the sets $P=\{x$ : $f(x)>0\}$ and $N=\{x: f(x) \leq 0\}$ form a Hahn decomposition of $X$ for $\nu$ and

$$
\nu^{+}(E)=\int_{E \cap P} f d \mu=\int_{E} f^{+} d \mu
$$

and

$$
\nu^{-}(E)=-\int_{E \cap N} f d \mu=\int_{E} f^{-} d \mu
$$

for each $E \in \mathcal{M}$.
6.39. Definition. The total variation of a signed measure $\nu$ is denoted as $\|\nu\|$ and is defined as

$$
\|\nu\|=\nu^{+}+\nu^{-}
$$

We conclude this section by examining alternate characterizations of absolutely continuous measures. We leave it as an exercise to prove that the following three conditions are equivalent:
(i) $\nu \ll \mu$
(ii) $\|\nu\| \ll \mu$
(iii) $\nu^{+} \ll \mu$ and $\nu^{-} \ll \mu$.
6.40. Theorem. Let $\nu$ be a finite signed measure and $\mu$ a positive measure on $(X, \mathcal{M})$. Then $\nu \ll \mu$ if and only if for every $\varepsilon>0$ there exists $\delta>0$ such that $|\nu(E)|<\varepsilon$ whenever $\mu(E)<\delta$.

Proof. Because of condition (ii) in (6.14) and the fact that $|\nu(E)| \leq\|\nu\|(E)$, we may assume that $\nu$ is a finite positive measure. Since the $\varepsilon, \delta$ condition is easily seen to imply that $\nu \ll \mu$, we will only prove the converse. Proceeding by contradiction, suppose then there exists $\varepsilon>0$ and a sequence of measurable sets $\left\{E_{k}\right\}$ such that $\mu\left(E_{k}\right)<2^{-k}$ and $\nu\left(E_{k}\right)>\varepsilon$ for all $k$. Set

$$
F_{m}=\bigcup_{k=m}^{\infty} E_{k}
$$

and

$$
F=\bigcap_{m=1}^{\infty} F_{m}
$$

Then, $\mu\left(F_{m}\right)<2^{1-m}$, so $\mu(F)=0$. But $\nu\left(F_{m}\right) \geq \varepsilon$ for each $m$ and, since $\nu$ is finite, we have

$$
\nu(F)=\lim _{m \rightarrow \infty} \nu\left(F_{m}\right) \geq \varepsilon
$$

thus reaching a contradiction.

## Exercises for Section 6.6

1. Prove property (iii) that follows (6.13).
2. Prove that the three conditions in (6.14) are equivalent.
3. Let $(X, \mathcal{M}, \mu)$ be a finite measure space, and let $f \in L^{1}(X, \mu)$. In particular, $f$ is $\mathcal{M}$-measurable. Suppose $\mathcal{M}_{0} \subset \mathcal{M}$ be a $\sigma$-algebra. Of course, $f$ may not be $\mathcal{M}_{0}$-measurable. However, prove that there is a unique $\mathcal{M}_{0}$-measurable function $f_{0}$ such that

$$
\int_{X} f g d \mu=\int_{X} f_{0} g d \tilde{\mu}
$$

for each $\mathcal{M}_{0}$-measurable $g$ for which the integrals are finite and $\tilde{\mu}=\mu\left\llcorner\mathcal{M}_{0}\right.$. Hint: Use the Radon-Nikodym Theorem.
4. Show that the total variation of the measure $\nu$ satisfies

$$
\|\nu\|(A)=\sup \left|\int_{A} f d \nu: f \in C_{c}(A),|f| \leq 1\right|
$$

for each open set $A$.
5. Suppose that $\mu$ and $\nu$ are $\sigma$-finite measures on $(X, \mathcal{M})$ such that $\mu \ll \nu$ and $\nu \ll \mu$. Prove that

$$
\frac{d \nu}{d \mu} \neq 0
$$

almost everywhere and

$$
\frac{d \mu}{d \nu}=1 / \frac{d \nu}{d \mu}
$$

6. Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a nondecreasing, continuously differentiable function and let $\lambda_{f}$ be the corresponding Lebesgue-Stieltjes measure, see Definition 4.29. Prove:
(a) $\lambda_{f} \ll \lambda$.
(b) $\frac{d \lambda_{f}}{d \lambda}=f^{\prime}$.
7. Let $(X, \mathcal{M}, \mu)$ be a finite measure space with $\mu(X)<\infty$. Let $\nu_{k}$ be a sequence of finite measures on $\mathcal{M}$ (that is, $\nu_{k}(X)<\infty$ for all $k$ ) with the property that they are uniformly absolutely continuous with respect to $\mu$; that is, for each $\varepsilon>0$, there exists $\delta>0$ and a positive integer $K$ such that $\nu_{k}(E)<\varepsilon$ for all $k \geq K$ and all $E \in \mathcal{M}$ for which $\mu(E)<\delta$. Assume that the limit

$$
\nu(E):=\lim _{k \rightarrow \infty} \nu_{k}(E), E \in \mathcal{M}
$$

exists. Prove that $\nu$ is a $\sigma$-finite measure on $\mathcal{M}$.
8. Let $(X, \mathcal{M}, \mu)$ a finite measure space and define a metric space $\widetilde{\mathcal{M}}$ as follows: for $A, B \in \mathcal{M}$, define

$$
\mathrm{d}(A, B):=\mu(A \Delta B) \quad \text { where } A \Delta B \text { denotes symmetric difference. }
$$

The space $\widetilde{\mathcal{M}}$ is defined as all sets in $\mathcal{M}$ where sets $A$ and $B$ are identified if $\mu(A \Delta B)=0$.
(a) Prove that $(\widetilde{\mathcal{M}}, \mathrm{d})$ is a complete metric space.
(b) Prove that $(\widetilde{\mathcal{M}}, \mathrm{d})$ is separable if and only if $L^{p}(X, \widetilde{\mathcal{M}}, \mu)$ is, $1 \leq p<\infty$.
9. Show that the space above is not compact when $X=[0,1], \mathcal{M}$ is the family of Borel sets on $[0,1]$, and $\mu$ is Lebesgue measure.
10. Let $\left\{\nu_{k}\right\}$ be a sequence of measures on the finite measure space $(X, \mathcal{M}, \mu)$ such that

- $\nu_{k}(X)<\infty$ for each $k$,
- the limit exists and is finite for each $E \in \mathcal{M}$

$$
\nu(E):=\lim _{k \rightarrow \infty} \nu_{k}(E),
$$

- $\nu_{k} \ll \mu$ for each $k$.
(a) Prove that each $\nu_{k}$ is well defined and continuous on the space $(\widetilde{\mathcal{M}}, \mathrm{d})$.
(b) For $\varepsilon>0$, let

$$
\mathcal{M}_{i, j}:=\left\{E \in \mathcal{M}:\left|\nu_{i}(E)-\nu_{j}(E)\right| \leq \frac{\varepsilon}{3}\right\}, i, j=1,2, \ldots
$$

and

$$
\mathcal{M}_{p}:=\bigcap_{i, j \geq p} \mathcal{M}_{i, j}, p=1,2, \ldots
$$

Prove that $\mathcal{M}_{p}$ is a closed set in $(\widetilde{\mathcal{M}}, \mathrm{d})$.
(c) Prove that there is some $q$ such that $\mathcal{M}_{q}$ contains an open set, call it $U$.
(d) Prove that the $\left\{\nu_{k}\right\}$ are uniformly absolutely continuous with respect to $\mu$, as in previous Problem 7.
Hint: Let $A$ be an interior point in $U$ and for $B \in \mathcal{M}$, write

$$
\nu_{k}(B)=\nu_{q}(B)+\left[\nu_{k}(B)-\nu_{q}(B)\right]
$$

and use the identity $\nu_{i}(B)=\nu_{i}(B \cap A)+\nu_{i}(B \backslash A), i=1,2 \ldots$ to estimate $\nu_{i}(B)$.
(e) Prove that $\nu$ is a finite measure.

### 6.7. The Radon-Nikodym Theorem

If $f$ is an integrable function on the measure space $(X, \mathcal{M}, \mu)$, then the signed measure

$$
\nu(E)=\int_{E} f d \mu
$$

defined for all $E \in \mathcal{M}$ is absolutely continuous with respect to $\mu$. The Radon-Nikodym theorem states that essentially every signed measure $\nu$, absolutely continuous with respect to $\mu$, is of this form. The proof of Theorem 6.41 below is due to A. Schep, [45].
6.41. Theorem. Suppose $(X, \mathcal{M}, \mu)$ is a finite measure space and $\nu$ is a measure on $(X, \mathcal{M})$ with the property

$$
\nu(E) \leq \mu(E)
$$

for each $E \in \mathcal{M}$. Then there is a measurable function $f: X \rightarrow[0,1]$ such that

$$
\begin{equation*}
\nu(E)=\int_{E} f d \mu, \quad \text { for each } \quad E \in \mathcal{M} \tag{6.15}
\end{equation*}
$$

More generally, if $g$ is a nonnegative measurable function on $X$, then

$$
\int_{X} g d \nu=\int_{X} g f d \mu
$$

Proof. Let

$$
H:=\left\{f: f \text { measurable, } 0 \leq f \leq 1, \int_{E} f d \mu \leq \nu(E) \text { for all } E \in \mathcal{M}\right\}
$$

and let

$$
M:=\sup \left\{\int_{X} f d \mu: f \in H\right\}
$$

Then, there exist functions $f_{k} \in H$ such that

$$
\int_{X} f_{k} d \mu>M-k^{-1}
$$

Observe that we may assume $0 \leq f_{1} \leq f_{2} \leq \ldots$ because if $f, g \in H$, then so is $\max \{f, g\}$ in view of the following:

$$
\begin{aligned}
\int_{E} \max \{f, g\} d \mu & =\int_{E \cap A} \max \{f, g\} d \mu+\int_{E \cap(X \backslash A)} \max \{f, g\} d \mu \\
& =\int_{E \cap A} f d \mu+\int_{X \backslash A} g d \mu \\
& \leq \nu(E \cap A)+\nu(E \cap(X \backslash A))=\nu(E)
\end{aligned}
$$

where $A:=\{x: f(x) \geq g(x)\}$. Therefore, since $\left\{f_{k}\right\}$ is an increasing sequence, the limit below exists:

$$
f_{\infty}(x):=\lim _{k \rightarrow \infty} f_{k}(x)
$$

Note that $f_{\infty}$ is a measurable function. Clearly $0 \leq f_{\infty} \leq 1$ and the Monotone Convergence theorem implies

$$
\int_{X} f_{\infty} d \mu=\lim _{k \rightarrow \infty} \int_{X} f_{k}=M
$$

for each $E \in \mathcal{M}$ and

$$
\int_{E} f_{k} d \mu \leq \nu(E)
$$

for each $k$ and $E \in \mathcal{M}$. So $f_{\infty} \in H$.
The proof of the theorem will be concluded by showing that (6.15) is satisfied by taking $f$ as $f_{\infty}$. For this purpose, assume by contradiction that

$$
\begin{equation*}
\int_{E} f_{\infty} d \mu<\nu(E) \tag{6.16}
\end{equation*}
$$

for some $E \in \mathcal{M}$. Let

$$
\begin{aligned}
& E_{0}=\left\{x \in E: f_{\infty}(x)<1\right\} \\
& E_{1}=\left\{x \in E: f_{\infty}(x)=1\right\}
\end{aligned}
$$

Then

$$
\begin{aligned}
\nu(E) & =\nu\left(E_{0}\right)+\nu\left(E_{1}\right) \\
& >\int_{E} f_{\infty} d \mu \\
& =\int_{E_{0}} f_{\infty} d \mu+\mu\left(E_{1}\right) \\
& \geq \int_{E_{0}} f_{\infty} d \mu+\nu\left(E_{1}\right)
\end{aligned}
$$

which implies

$$
\nu\left(E_{0}\right)>\int_{E_{0}} f_{\infty} d \mu
$$

Let $\varepsilon^{*}>0$ be such that

$$
\begin{equation*}
\int_{E_{0}} f_{\infty}+\varepsilon^{*} \chi_{E_{0}} d \mu<\nu\left(E_{0}\right) \tag{6.17}
\end{equation*}
$$

and let $F_{k}:=\left\{x \in E_{0}: f_{\infty}(x)<1-1 / k\right\}$. Observe that $F_{1} \subset F_{2} \subset \ldots$ and since $f_{\infty}<1$ on $E_{0}$, we have $\cup_{k=1}^{\infty} F_{k}=E_{0}$ and therefore that $\nu\left(F_{k}\right) \uparrow \nu\left(E_{0}\right)$. Furthermore, it follows from (6.17) that

$$
\int_{E_{0}} f_{\infty}+\varepsilon^{*} \chi_{E_{0}} d \mu<\nu\left(E_{0}\right)-\eta
$$

for some $\eta>0$. Therefore, since $\left[f_{\infty}+\varepsilon^{*} \chi_{F_{k}}\right] \chi_{F_{k}} \uparrow\left[f_{\infty}+\varepsilon^{*} \chi_{E_{0}}\right] \chi_{E_{0}}$, there exists $k^{*}$ such that

$$
\begin{aligned}
& \int_{F_{k}} f_{\infty}+\varepsilon^{*} \chi_{F_{k}} d \mu \\
& =\int_{X}\left[f_{\infty}+\varepsilon^{*} \chi_{F_{k}}\right] \chi_{F_{k}} d \mu \\
& \rightarrow \int_{E_{0}} f_{\infty}+\varepsilon^{*} \chi_{E_{0}} d \mu \quad \text { by Monotone Convergence Theorem } \\
& <\nu\left(E_{0}\right)-\eta \\
& <\nu\left(F_{k}\right)
\end{aligned} \quad \text { for all } k \geq k^{*} .
$$

For all such $k \geq k^{*}$ and $\varepsilon:=\min \left(\varepsilon^{*}, 1 / k\right)$, we claim that

$$
\begin{equation*}
f_{\infty}+\varepsilon \chi_{F_{k}} \in H \tag{6.18}
\end{equation*}
$$

The validity of this claim would imply that $\int f_{\infty}+\varepsilon \chi_{F_{k}} d \mu=M+\varepsilon \mu\left(F_{k}\right)>$ $M$, contradicting the definition of $M$ which would mean that our contradiction hypothesis, (6.16), is false, thus establishing our theorem.

So, to finish the proof, it suffices to prove (6.18) for any $k \geq k^{*}$, which will remain fixed throughout the remainder of the proof. For this, first note that $0 \leq$ $f_{\infty}+\varepsilon \chi_{F_{k}} \leq 1$. To show that

$$
\int_{E} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \leq \nu(E) \quad \text { for all } E \in \mathcal{M}
$$

we proceed by contradiction; if not, there would exist a measurable set $G \subset X$ such that

$$
\begin{aligned}
\nu\left(G \backslash F_{k}\right) & +\int_{G \cap F_{k}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \\
& \geq \int_{G \backslash F_{k}} f_{\infty} d \mu+\int_{G \cap F_{k}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \\
& =\int_{G \backslash F_{k}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu+\int_{G \cap F_{k}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \\
& =\int_{G} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu>\nu(G) .
\end{aligned}
$$

This implies

$$
\begin{equation*}
\int_{G \cap F_{k}}\left(f_{\infty}+\varepsilon \chi_{F_{k}}\right) d \mu>\nu(G)-\nu\left(G \backslash F_{k}\right)=\nu\left(G \cap F_{k}\right) \tag{6.19}
\end{equation*}
$$

Hence, we may assume $G \subset F_{k}$.
Let $\mathcal{F}_{1}$ be the collection of all measurable sets $G \subset F_{k}$ such that (6.19) holds. Define $\alpha_{1}:=\sup \left\{\mu(G): G \in \mathcal{F}_{1}\right\}$ and let $G_{1} \in \mathcal{F}_{1}$ be such that $\mu\left(G_{1}\right)>\alpha_{1}-1$. Similarly, let $\mathcal{F}_{2}$ be the collection of all measurable sets $G \subset F_{k} \backslash G_{1}$ such that
$\left.{ }^{*}\right)$ holds for $G$. Define $\alpha_{2}:=\sup \left\{\mu(G): G \in \mathcal{F}_{2}\right\}$ and let $G_{2} \in \mathcal{F}_{2}$ be such that $\mu\left(G_{2}\right)>\alpha_{2}-\frac{1}{2^{2}}$. Proceeding inductively, we obtain a decreasing sequence $\alpha_{k}$ and disjoint measurable sets $G_{j}$ where $\mu\left(G_{j}\right)>\alpha_{j}-\frac{1}{j^{2}}$. Observe that $\alpha_{j} \rightarrow 0$, for if $\alpha_{j} \rightarrow a>0$, then $\mu\left(G_{j}\right) \downarrow a$. Since $\sum_{j=1}^{\infty} \frac{1}{j^{2}}<\infty$ this would imply

$$
\mu\left(\bigcup G_{j}\right)=\sum_{j=1}^{\infty} \mu\left(G_{j}\right)>\sum_{j=1}^{\infty}\left(\alpha_{j}-\frac{1}{j^{2}}\right)=\infty
$$

contrary to the finiteness of $\mu$. This implies

$$
\mu\left(F_{k} \backslash \bigcup_{j=1}^{\infty} G_{j}\right)=0
$$

for if not, there would be two possibilities:
(i) there would exist a set $T \subset F_{k} \backslash \bigcup_{j=1}^{\infty} G_{j}$ of positive $\mu$ measure for which (6.19) would hold with $T$ replacing $G$. Since $\mu(T)>0$, there would exist $\alpha_{k}$ such that

$$
\alpha_{k}<\mu(T)
$$

and since

$$
T \subset F_{k} \backslash \bigcup_{j=1}^{\infty} G_{j} \subset F_{k} \backslash \bigcup_{i=1}^{j-1} G_{j}
$$

this would contradict the definition of $\alpha_{j}$. Hence,

$$
\begin{aligned}
\nu\left(F_{k}\right) & >\int_{F_{k}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \\
& =\sum_{j} \int_{G_{j}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \\
& >\sum_{j} \nu\left(G_{j}\right)=\nu\left(F_{k}\right)
\end{aligned}
$$

which is impossible and therefore (i) cannot occur.
(ii) If there were no set $T$ as in (i), then $F_{k} \backslash \bigcup_{j}^{\infty} G_{j}$ could not satisfy (6.18) and thus

$$
\begin{equation*}
\int_{F_{k} \backslash \cup_{j} G_{j=1}^{\infty}} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \leq \nu\left(F_{k} \backslash \cup_{j=1}^{\infty} G_{j}\right) \tag{6.20}
\end{equation*}
$$

With $S:=F_{k} \backslash \bigcup_{j=1}^{\infty} G_{j}$ we have

$$
\int_{S} f_{\infty}+\varepsilon \chi_{F_{k}} d \mu \leq \nu(S)
$$

Since

$$
\int_{G_{j}} f_{\infty}+\varepsilon \chi_{F_{k}}>\nu\left(G_{j}\right)
$$

for each $j \in \mathbb{N}$, it follows that $F_{k} \backslash S$ must also satisfy (6.18), which contradicts (6.20). Hence, both (i) and (ii) do not occur and thus we conclude that

$$
\mu\left(F_{k} \backslash \bigcup_{j=1}^{\infty} G_{j}\right)=0
$$

as desired.
6.42. Notation. The function $f$ in (6.15) (and also in (6.22) below) is called the Radon-Nikodym derivative of $\nu$ with respect to $\mu$ and is denoted by

$$
f:=\frac{d \nu}{d \mu}
$$

The previous theorem yields the notationally convenient result

$$
\begin{equation*}
\int_{X} g d \nu=\int_{X} g \frac{d \nu}{d \mu} d \mu \tag{6.21}
\end{equation*}
$$

6.43. Theorem (Radon-Nikodym). If $(X, \mathcal{M}, \mu)$ is a $\sigma$-finite measure space and $\nu$ is a $\sigma$-finite signed measure on $\mathcal{M}$ that is absolutely continuous with respect to $\mu$, then there exists a measurable function $f$ such that either $f^{+}$or $f^{-}$is integrable and

$$
\begin{equation*}
\nu(E)=\int_{E} f d \mu \tag{6.22}
\end{equation*}
$$

for each $E \in \mathcal{M}$.
Proof. We first assume, temporarily, that $\mu$ and $\nu$ are finite measures. Referring to Theorem 6.41, there exist Radon-Nikodym derivatives

$$
f_{\nu}:=\frac{d \nu}{d(\nu+\mu)} \quad \text { and } \quad f_{\mu}:=\frac{d \mu}{d(\nu+\mu)}
$$

Define $A=X \cap\left\{f_{\mu}(x)>0\right\}$ and $B=X \cap\left\{f_{\mu}(x)=0\right\}$. Then

$$
\mu(B)=\int_{B} f_{\mu} d(\nu+\mu)=0
$$

and therefore, $\nu(B)=0$ since $\nu \ll \mu$. Now define

$$
f(x)= \begin{cases}\frac{f_{\nu}(x)}{f_{\mu}(x)} & \text { if } x \in A \\ 0 & \text { if } x \in B\end{cases}
$$

If $E$ is a measurable subset of $A$, then

$$
\nu(E)=\int_{E} f_{\nu} d(\nu+\mu)=\int_{E} f \cdot f_{\mu} d(\nu+\mu)=\int_{E} f d \mu
$$

by (6.21). Since both $\nu$ and $\mu$ are 0 on $B$, we have

$$
\nu(E)=\int_{E} f d \mu
$$

for all measurable $E$.

Next, consider the case where $\mu$ and $\nu$ are $\sigma$-finite measures. There is a sequence of disjoint measurable sets $\left\{X_{k}\right\}_{k=1}^{\infty}$ such that $X=\cup_{k=1}^{\infty} X_{k}$ and both $\mu\left(X_{k}\right)$ and $\nu\left(X_{k}\right)$ are finite for each $k$. Set $\mu_{k}:=\mu\left\llcorner X_{k}, \nu_{k}:=\nu\left\llcorner X_{k}\right.\right.$ for $k=1,2, \ldots$. Clearly $\mu_{k}$ and $\nu_{k}$ are finite measures on $\mathcal{M}$ and $\nu_{k} \ll \mu_{k}$. Thus there exist nonnegative measurable functions $f_{k}$ such that for any $E \in \mathcal{M}$

$$
\nu\left(E \cap X_{k}\right)=\nu_{k}(E)=\int_{E} f_{k} d \mu_{k}=\int_{E \cap X_{k}} f_{k} d \mu
$$

It is clear that we may assume $f_{k}=0$ on $X-X_{k}$. Set $f:=\sum_{k=1}^{\infty} f_{k}$. Then for any $E \in \mathcal{M}$

$$
\nu(E)=\sum_{k=1}^{\infty} \nu\left(E \cap X_{k}\right)=\sum_{k=1}^{\infty} \int_{E \cap X_{k}} f d \mu=\int_{E} f d \mu .
$$

Finally suppose that $\nu$ is a signed measure and let $\nu=\nu^{+}-\nu^{-}$be the Jordan decomposition of $\nu$. Since the measures are mutually singular there is a measurable set $P$ such that $\nu^{+}(X-P)=0=\nu^{-}(P)$. For any $E \in \mathcal{M}$ such that $\mu(E)=0$,

$$
\begin{aligned}
& \nu^{+}(E)=\nu(E \cap P)=0 \\
& \nu^{-}(E)=-\nu(E-P)=0
\end{aligned}
$$

since $\mu(E \cap P)+\mu(E-P)=\mu(E)=0$. Thus $\nu^{+}$and $\nu^{-}$are absolutely continuous with respect to $\mu$ and consequently, there exist nonnegative measurable functions $f^{+}$and $f^{-}$such that

$$
\nu^{ \pm}(E)=\int_{E} f^{ \pm} d \mu
$$

for each $E \in \mathcal{M}$. Since at least one of the measures $\nu^{ \pm}$is finite it follows that at least one of the functions $f^{ \pm}$is $\mu$-integrable. Set $f=f^{+}-f^{-}$. In view of Theorem 6.9 , p. 157,

$$
\nu(E)=\int_{E} f^{+} d \mu-\int_{E} f^{-} d \mu=\int_{E} f d \mu
$$

for each $E \in \mathcal{M}$.

An immediate consequence of this result is the following.
6.44. Theorem (Lebesgue Decomposition). Let $\mu$ and $\nu$ be $\sigma$-finite measures defined on the measure space $(X, \mathcal{M})$. Then there is a decomposition of $\nu$ such that $\nu=\nu_{0}+\nu_{1}$ where $\nu_{0} \perp \mu$ and $\nu_{1} \ll \mu$. The measures $\nu_{0}$ and $\nu_{1}$ are unique.

Proof. We employ the same device as in the proof of the preceding theorem by considering the Radon-Nikodym derivatives

$$
f_{\nu}:=\frac{d \nu}{d(\nu+\mu)} \quad \text { and } \quad f_{\mu}:=\frac{d \mu}{d(\nu+\mu)}
$$

Define $A=X \cap\left\{f_{\mu}(x)>0\right\}$ and $B=X \cap\left\{f_{\mu}(x)=0\right\}$. Then $X$ is the disjoint union of $A$ and $B$. With $\gamma:=\mu+\nu$ we will show that the measures

$$
\nu_{0}(E):=\nu(E \cap B) \quad \text { and } \quad \nu_{1}(E):=\nu(E \cap A)=\int_{E \cap A} f_{\nu} d \gamma
$$

provide our desired decomposition. First, note that $\nu=\nu_{0}+\nu_{1}$. Next, we have $\nu_{0}(A)=0$ and so $\nu_{0} \perp \mu$. Finally, to show that $\nu_{1} \ll \mu$, consider $E$ with $\mu(E)=0$. Then

$$
0=\mu(E)=\int_{E} f_{\mu} d \gamma=\int_{A \cap E} f_{\mu} d \lambda
$$

Thus, $f_{\mu}=0 \gamma$-a.e. on $E$. Then, since $f_{\mu}>0$ on $A$, we must have $\gamma(A \cap E)=0$. This implies $\nu(A \cap E)=0$ and therefore $\nu_{1}(E)=0$, which establishes $\nu_{1} \ll \mu$. The proof of uniqueness is left as an exercise.

## Exercises for Section 6.7

1. Prove the uniqueness assertion in Theorem 6.44.

### 6.8. The Dual of $L^{p}$

Using the Radon-Nikodym Theorem we completely characterize the continuous linear mappings of $L^{p}(X)$ into $\mathbb{R}$.
6.45. Definitions. Let $(X, \mathcal{M}, \mu)$ be a measure space. A linear functional on $L^{p}(X)=L^{p}(X, \mathcal{M}, \mu)$ is a real-valued linear function on $L^{p}(X)$, i.e., a function $F: L^{p}(X) \rightarrow \mathbb{R}$ such that

$$
F(a f+b g)=a F(f)+b F(g)
$$

whenever $f, g \in L^{p}(X)$ and $a, b \in \mathbb{R}$. Set

$$
\|F\| \equiv \sup \left\{|F(f)|: f \in L^{p}(X),\|f\|_{p} \leq 1\right\}
$$

A linear functional $F$ on $L^{p}(X)$ is said to be bounded if $\|F\|<\infty$.
6.46. Theorem. A linear functional on $L^{p}(X)$ is bounded if and only if it is continuous with respect to (the metric induced by) the norm $\|\cdot\|_{p}$.

Proof. Let $F$ be a linear functional on $L^{p}(X)$.
If $F$ is bounded, then $\|F\|<\infty$ and if $0 \neq f \in L^{p}(X)$ then

$$
\left|F\left(\frac{f}{\|f\|_{p}}\right)\right| \leq\|F\|
$$

i.e.,

$$
|F(f)| \leq\|F\|\|f\|_{p}
$$

whenever $f \in L^{p}(X)$. In particular for any $f, g \in L^{p}(X)$

$$
|F(f-g)| \leq\|F\|\|f-g\|_{p}
$$

and hence $F$ is uniformly continuous on $L^{p}(X)$.
On the other hand if $F$ is continuous at 0 , then there exists a $\delta>0$ such that

$$
|F(f)| \leq 1
$$

whenever $\|f\|_{p} \leq \delta$. Thus if $f \in L^{p}(X)$ with $\|f\|_{p}>0$, then

$$
|F(f)|=\left|\frac{\|f\|_{p}}{\delta} F\left(\frac{\delta}{\|f\|_{p}} f\right)\right| \leq \frac{1}{\delta}\|f\|_{p}
$$

whence $\|F\| \leq \frac{1}{\delta}$.
6.47. Theorem. If $1 \leq p \leq \infty, \frac{1}{p}+\frac{1}{p^{\prime}}=1$ and $g \in L^{p^{\prime}}(X)$, then

$$
F(f)=\int f g d \mu
$$

defines a bounded linear functional on $L^{p}(X)$ with

$$
\|F\|=\|g\|_{p^{\prime}}
$$

Proof. That $F$ is a bounded linear functional on $L^{p}(X)$ follows immediately from Hölder's inequality and the elementary properties of the integral. The rest of the assertion follows from Theorem 6.26 since

$$
\|F\|=\sup \left\{\int f g d \mu:\|f\|_{p} \leq 1\right\}=\|g\|_{p^{\prime}}
$$

Note that while the hypotheses of Theorem 6.26 , include a $\sigma$-finiteness condition, that assumption is not needed to establish (6.12) if the function is integrable. That's the situation we have here since it is assumed that $g \in L^{p^{\prime}}$.

The next theorem shows that all bounded linear functionals on $L^{p}(X)(1 \leq$ $p<\infty)$ are of this form.
6.48. TheOrem. If $1<p<\infty$ and $F$ is a bounded linear functional on $L^{p}(X)$, then there is a $g \in L^{p^{\prime}}(X),\left(\frac{1}{p}+\frac{1}{p^{\prime}}=1\right)$ such that

$$
\begin{equation*}
F(f)=\int f g d \mu \tag{6.23}
\end{equation*}
$$

for all $f \in L^{p}(X)$. Moreover $\|g\|_{p^{\prime}}=\|F\|$ and the function $g$ is unique in the sense that if (6.23) holds with $\tilde{g} \in L^{p^{\prime}}(X)$, then $\tilde{g}=g \mu$-a.e.. If $p=1$, the same conclusion holds under the additional assumption that $\mu$ is $\sigma$-finite.

Proof. Assume first $\mu(X)<\infty$. Note that our assumption imply that $\chi_{E} \in$ $L^{p}(X)$ whenever $E \in \mathcal{M}$. Set

$$
\nu(E)=F\left(\chi_{E}\right)
$$

for $E \in \mathcal{M}$. Suppose $\left\{E_{k}\right\}_{k=1}^{\infty}$ is a sequence of disjoint measurable sets and let $E:=\bigcup_{k=1}^{\infty} E_{k}$. Then for any positive integer $N$

$$
\begin{aligned}
\left|\nu(E)-\sum_{k=1}^{N} \nu\left(E_{k}\right)\right| & =\left|F\left(\chi_{E}-\sum_{k=1}^{N} \chi_{E_{k}}\right)\right| \\
& =\left|F\left(\sum_{k=N+1}^{\infty} \chi_{E_{k}}\right)\right| \\
& \leq\|F\|\left(\mu\left(\bigcup_{k=N+1}^{\infty} E_{k}\right)\right)^{\frac{1}{p}}
\end{aligned}
$$

and $\mu\left(\bigcup_{k=N+1}^{\infty} E_{k}\right)=\sum_{k=N+1}^{\infty} \mu\left(E_{k}\right) \rightarrow 0$ as $N \rightarrow \infty$ since

$$
\mu(E)=\sum_{k=1}^{\infty} \mu\left(E_{k}\right)<\infty
$$

Thus,

$$
\nu(E)=\sum_{k=1}^{\infty} \nu\left(E_{k}\right)
$$

and since the same result holds for any rearrangement of the sequence $\left\{E_{k}\right\}_{k=1}^{\infty}$, the series converges absolutely. It follows that $\nu$ is a signed measure and since $|\nu(E)| \leq\|F\|(\mu(E))^{\frac{1}{p}}$ we see that $\nu \ll \mu$. By the Radon-Nikodym Theorem there is a $g \in L^{1}(X)$ such that

$$
F\left(\chi_{E}\right)=\nu(E)=\int \chi_{E} g d \mu
$$

for each $E \in \mathcal{M}$. From the linearity of both $F$ and the integral, it is clear that

$$
\begin{equation*}
F(f)=\int f g d \mu \tag{6.24}
\end{equation*}
$$

whenever $f$ is a simple function.

Step 1: Assume $\mu(X)<\infty$ and $p=1$.
We proceed to show that $g \in L^{\infty}(X)$. Assume that $\left\|g^{+}\right\|_{\infty}>\|F\|$. Let $M>0$ such that $\left\|g^{+}\right\|_{\infty}>M>\|F\|$ and set $E_{M}:=\{x: g(x)>M\}$. We have $\mu\left(E_{M}\right)>0$, since otherwise we would have $\left\|g^{+}\right\|_{\infty} \leq M$. Then

$$
M \mu\left(E_{M}\right) \leq \int \chi_{E_{M}} g d \mu=F\left(\chi_{E_{M}}\right) \leq\|F\| \mu\left(E_{M}\right)
$$

Thus $\mu\left(E_{M}\right)>0$ yields $M \leq\|F\|$, which is a contradiction. We conclude that $\left\|g^{+}\right\|_{\infty} \leq\|F\|$. Similarly $\left\|g^{-}\right\|_{\infty} \leq\|F\|$ and hence $\|g\|_{\infty} \leq\|F\|$.

If $f$ is an arbitrary function in $L^{1}$, then we know by Theorem 5.27, that there exist simple functions $f_{k}$ with $\left|f_{k}\right| \leq|f|$ such that $\left\{f_{k}\right\} \rightarrow f$ pointwise and $\left\|f_{k}-f\right\|_{1} \rightarrow 0$. Therefore, by Lebesgue's Dominated Convergence theorem,

$$
\left|F\left(f_{k}\right)-F(f)\right|=\left|F\left(f_{k}-f\right)\right| \leq\|F\|\left\|f-f_{k}\right\|_{1} \rightarrow 0
$$

and

$$
\begin{aligned}
\left|F(f)-\int f g d \mu\right| & \leq\left|F\left(f-f_{k}\right)\right|+\left|\int\left(f_{k}-f\right) g d \mu\right| \\
& \leq\|F\|\left\|f-f_{k}\right\|_{1}+\left\|f-f_{k}\right\|_{1}\|g\|_{\infty} \\
& \leq 2\|F\|\left\|f-f_{k}\right\|_{1}
\end{aligned}
$$

and thus we have our desired result when $p=1$ and $\mu(X)<\infty$. By Theorem 6.26, we have $\|g\|_{\infty}=\|F\|$.

Step 2. Assume $\mu(X)<\infty$ and $1<p<\infty$.
Let $\left\{h_{k}\right\}_{k=1}^{\infty}$ be an increasing sequence of nonnegative simple functions such that $\lim _{k \rightarrow \infty} h_{k}=|g|$. Set $g_{k}=h_{k}^{p^{\prime}-1} \operatorname{sign}(g)$. Then

$$
\begin{align*}
\left\|h_{k}\right\|_{p^{\prime}}^{p^{\prime}} & =\int\left|h_{k}\right|^{p^{\prime}} d \mu \leq \int g_{k} g d \mu=F\left(g_{k}\right) \\
& \leq\|F\|\left\|g_{k}\right\|_{p}=\|F\|\left(\int h_{k}^{p^{\prime}} d \mu\right)^{\frac{1}{p}}=\|F\|\left\|h_{k}\right\|_{p^{\prime}}^{\frac{p^{\prime}}{p}} . \tag{6.25}
\end{align*}
$$

We wish to conclude that $g \in L^{p^{\prime}}(X)$. For this we may assume that $\|g\|_{p^{\prime}}>0$ and hence that $\left\|h_{k}\right\|_{p^{\prime}}>0$ for large $k$. It then follows from (6.25) that $\left\|h_{k}\right\|_{p^{\prime}} \leq\|F\|$ for all $k$ and thus by Fatou's Lemma we have

$$
\|g\|_{p^{\prime}} \leq \liminf _{k \rightarrow \infty}\left\|h_{k}\right\|_{p^{\prime}} \leq\|F\|
$$

which shows that $g \in L^{p^{\prime}}(X), 1<p<\infty$.
Now let $f \in L^{p}(X)$ and let $\left\{f_{k}\right\}$ be a sequence of simple functions such that $\left\|f-f_{k}\right\|_{p} \rightarrow 0$ as $k \rightarrow \infty$ (Exercise 1, Section 6.5). Then

$$
\begin{aligned}
\left|F(f)-\int f g d \mu\right| & \leq\left|F\left(f-f_{k}\right)\right|+\left|\int\left(f_{k}-f\right) g d \mu\right| \\
& \leq\|F\|\left\|f-f_{k}\right\|_{p}+\left\|f-f_{k}\right\|_{p}\|g\|_{p^{\prime}} \\
& \leq 2\|F\|\left\|f-f_{k}\right\|_{p}
\end{aligned}
$$

for all $k$; whence,

$$
F(f)=\int f g d \mu
$$

for all $f \in L^{p}(X)$. By Theorem 6.26 , we have $\|g\|_{p^{\prime}}=\|F\|$. Thus, using step 1 also, we conclude that the proof is complete under the assumptions that $\mu(X)<\infty$, $1 \leq p<\infty$.

Step 3. Assume $\mu$ is $\sigma$-finite and $1 \leq p<\infty$.
Suppose $Y \in \mathcal{M}$ is $\sigma$-finite. Let $\left\{Y_{k}\right\}$ be an increasing sequence of measurable sets such that $\mu\left(Y_{k}\right)<\infty$ for each $k$ and such that $Y=\bigcup_{k=1}^{\infty} Y_{k}$. Then, from Steps 1 and 2 above (see also Exercise 7, Section 6.1), for each $k$ there is a measurable function $g_{k}$ such that $\left\|g_{k} \chi_{Y_{k}}\right\|_{p^{\prime}} \leq\|F\|$ and

$$
F\left(f \chi_{Y_{k}}\right)=\int f \chi_{Y_{k}} g_{k} d \mu
$$

for each $f \in L^{p}(X)$. We may assume $g_{k}=0$ on $Y-Y_{k}$. If $k<m$ then

$$
F\left(f \chi_{Y_{k}}\right)=\int f g_{m} \chi_{Y_{k}} d \mu
$$

for each $f \in L^{p}(X)$. Thus

$$
\int f\left(g_{k}-g_{m} \chi_{Y_{k}}\right) d \mu=0
$$

for each $f \in L^{p}(X)$. By Theorem 6.26, this implies that $g_{k}=g_{m} \mu$-a.e. on $Y_{k}$. Thus $\left\{g_{k}\right\}$ converges $\mu$-a.e. to a measurable function $g$ and by Fatou's Lemma

$$
\begin{equation*}
\|g\|_{p^{\prime}} \leq \liminf _{k \rightarrow \infty}\left\|g_{k}\right\|_{p^{\prime}} \leq\|F\| \tag{6.26}
\end{equation*}
$$

which shows that $g \in L^{p^{\prime}}, 1<p^{\prime}<\infty$. For $p^{\prime}=\infty$, we also have $\|g\|_{\infty} \leq\|F\|$.
Fix $f \in L^{p}(X)$ and set $f_{k}:=f \chi_{Y_{k}}$. Then $f_{k}$ converges to $f \chi_{Y}$ and $\left|f-f_{k}\right| \leq$ $2|f|$. By the Dominated Convergence Theorem $\left\|\left(f \chi_{Y}-f_{k}\right)\right\|_{p} \rightarrow 0$ as $k \rightarrow \infty$. Thus, since $g_{k}=g \mu$-a.e. on $Y_{k}$,

$$
\begin{aligned}
\left|F\left(f \chi_{Y}\right)-\int f g d \mu\right| & \leq\left|F\left(f \chi_{Y}-f_{k}\right)\right|+\int\left|f-f_{k}\right||g| d \mu \\
& \leq 2\|F\|\left\|f \chi_{Y}-f_{k}\right\|_{p}
\end{aligned}
$$

and therefore

$$
\begin{equation*}
F\left(f \chi_{Y}\right)=\int f g d \mu \tag{6.27}
\end{equation*}
$$

for each $f \in L^{p}(X)$.
If $\mu$ is $\sigma$-finite, we may set $Y=X$ and deduce that

$$
F(f)=\int f g d \mu
$$

whenever $f \in L^{p}(X)$. Thus, in view of Theorem 6.47,

$$
\|F\|=\sup \left\{|F(f)|:\|f\|_{p}=1\right\}=\|g\|_{p^{\prime}}
$$

Step 4. Assume $\mu$ is not $\sigma$-finite and $1<p<\infty$.

When $1<p<\infty$ and $\mu$ is not assumed to be $\sigma$-finite, we will conclude the proof by making a judicious choice for $Y$ in (6.27) so that

$$
\begin{equation*}
F(f)=F\left(f \chi_{Y}\right) \quad \text { for each } f \in L^{p}(X) \text { and }\|g\|_{p^{\prime}}=\|F\| \tag{6.28}
\end{equation*}
$$

For each positive integer $k$ there is $h_{k} \in L^{p}(X)$ such that $\left\|h_{k}\right\|_{p} \leq 1$ and

$$
\|F\|-\frac{1}{k}<\left|F\left(h_{k}\right)\right| .
$$

Set

$$
Y=\bigcup_{k=1}^{\infty}\left\{x: h_{k}(x) \neq 0\right\}
$$

Then $Y$ is a measurable, $\sigma$-finite subset of $X$ and thus, by Step 3, there is a $g \in L^{p^{\prime}}(X)$ such that $g=0 \mu$-a.e. on $X-Y$ and

$$
\begin{equation*}
F\left(f \chi_{Y}\right)=\int f g d \mu \tag{6.29}
\end{equation*}
$$

for each $f \in L^{p}(X)$. Since for each $k$

$$
\|F\|-\frac{1}{k}<F\left(h_{k}\right)=\int h_{k} g d \mu \leq\|g\|_{p^{\prime}}
$$

we see that $\|g\|_{p^{\prime}} \geq\|F\|$. On the other hand, appealing to Theorem 6.26 again,

$$
\|F\| \geq \sup _{\|f\|_{p} \leq 1} F\left(f \chi_{Y}\right)=\sup _{\|f\|_{p} \leq 1} \int f g d \mu=\|g\|_{p^{\prime}}
$$

which establishes the second part of (6.28),
To establish the first part of (6.28), with the help of (6.29), it suffices to show that $F(f)=F\left(f \chi_{Y}\right)$ for each $f \in L^{p}(X)$. By contradiction, suppose there is a function $f_{0} \in L^{p}(X)$ such that $F\left(f_{0}\right) \neq F\left(f_{0} \chi_{Y}\right)$. Set $Y_{0}=\left\{x: f_{0}(x) \neq 0\right\}-Y$. Then, since $Y_{0}$ is $\sigma$-finite, there is a $g_{0} \in L^{p^{\prime}}(X)$ such that $g_{0}=0 \mu$-a.e. on $X-Y_{0}$ and

$$
F\left(f \chi_{Y_{0}}\right)=\int f g_{0} d \mu
$$

for each $f \in L^{p}(X)$. Since $g$ and $g_{0}$ are non-zero on disjoint sets, note that

$$
\left\|g+g_{0}\right\|_{p^{\prime}}^{p^{\prime}}=\|g\|_{p^{\prime}}^{p^{\prime}}+\left\|g_{0}\right\|_{p^{\prime}}^{p^{\prime}}
$$

and that $\left\|g_{0}\right\|_{p^{\prime}}>0$ since

$$
\int f_{0} g_{0} d \mu=F\left(f_{0}\left[1-\chi_{Y}\right]\right)=F\left(f_{0}\right)-F\left(f_{0} \chi_{Y}\right) \neq 0
$$

Moreover, since

$$
F\left(f \chi_{Y \cup Y_{0}}\right)=\int f\left(g+g_{0}\right) d \mu
$$

for each $f \in L^{p}(X)$, we see that $\left\|g+g_{0}\right\|_{p^{\prime}} \leq\|F\|$. Thus

$$
\begin{aligned}
\|F\|^{p^{\prime}} & =\|g\|_{p^{\prime}}^{p^{\prime}} \\
& <\|g\|_{p^{\prime}}^{p^{\prime}}+\left\|g_{0}\right\|_{p^{\prime}}^{p^{\prime}} \\
& =\left\|g+g_{0}\right\|_{p^{\prime}}^{p^{\prime}} \\
& \leq\|F\|^{p^{\prime}} .
\end{aligned}
$$

This contradiction implies that

$$
F(f)=\int f g d \mu
$$

for each $f \in L^{p}(X)$, which establishes the first part of (6.28), as desired.
Step 5. Uniqueness of $g$.
If $\tilde{g} \in L^{p^{\prime}}(X)$ is such that

$$
\int f(g-\tilde{g}) d \mu=0
$$

for all $f \in L^{p}(X)$, then by Theorem $6.26\|g-\tilde{g}\|_{p^{\prime}}=0$ and thus $g=\tilde{g} \mu$-a.e. thus establishing the uniqueness of $g$.

## Exercises for Section 6.8

1. Suppose $f$ is a nonnegative measurable function. Set

$$
E_{t}=\{x: f(x)>t\}
$$

and

$$
g(t)=-\mu\left(E_{t}\right)
$$

for $t \in \mathbb{R}$. Show that

$$
\int f d \mu=\int_{0}^{\infty} t d \lambda_{g}(t)
$$

where $\lambda_{g}$ is the Lebesgue-Stieltjes measure induced by $g$ as in Section 4.6.
2. Let $X$ be a well-ordered set (with ordering denoted by $<$ ) that is a representative of the ordinal number $\Omega$ and let $\mathcal{M}$ be the $\sigma$-algebra consisting of those sets $E$ with the property that either $E$ or its complement is at most countable. Let $\mu$ be the measure defined on $E \in \mathcal{M}$ as $\mu(E)=0$ if $E$ is at most countable and $\mu(E)=1$ otherwise. Let $Y:=X, \nu:=\mu$ and
(a) Show that if $A:=\{(x, y) \in X \times Y: x<y\}$, then $A_{x}$ and $A_{y}$ are measurable for all $x$ and $y$.
(b) Show that both

$$
\int_{Y}\left(\int_{X} \chi_{A}(x, y) d \mu(x)\right) d \nu(y)
$$

and

$$
\int_{X}\left(\int_{Y} \chi_{A}(x, y) d \nu(y)\right) d \mu(x)
$$

exist.
(c) Show that

$$
\int_{Y}\left(\int_{X} \chi_{A}(x, y) d \mu(x)\right) d \nu(y) \neq \int_{X}\left(\int_{Y} \chi_{A}(x, y) d \nu(y)\right) d \mu(x)
$$

(d) Why doesn't Fubini's Theorem apply in this example?
3. Let $f$ and $g$ be integrable functions on a measure space $(X, \mathcal{M}, \mu)$ with the property that

$$
\mu[\{f>t\} \Delta\{g>t\}]=0
$$

for $\lambda$-a.e. $t$. Prove that $f=g \mu$-a.e.
4. Let $f$ be a Lebesgue measurable function on $[0,1]$ and let $Q:=[0,1] \times[0,1]$.
(a) Show that $F(x, y):=f(x)-f(y)$ is measurable with respect to Lebesgue measure in $\mathbb{R}^{2}$.
(b) If $F \in L^{1}(Q)$, show that $f \in L^{1}([0,1])$.

### 6.9. Product Measures and Fubini's Theorem

In this section we introduce product measures and prove Fubini's theorem, which generalizes the notion of iterated integration of Riemannian calculus.

Let $\left(X, \mathcal{M}_{X}, \mu\right)$ and $\left(Y, \mathcal{M}_{Y}, \nu\right)$ be two complete measure spaces. In order to define the product of $\mu$ and $\nu$ we first define an outer measure on $X \times Y$ in terms of $\mu$ and $\nu$.
6.49. Definition. For each $S \subset X \times Y$ set

$$
\begin{equation*}
\varphi(S)=\inf \left\{\sum_{j=1}^{\infty} \mu\left(A_{j}\right) \nu\left(B_{j}\right)\right\} \tag{6.30}
\end{equation*}
$$

where the infimum is taken over all sequences $\left\{A_{j} \times B_{j}\right\}_{j=1}^{\infty}$ such that $A_{j} \in \mathcal{M}_{X}$, $B_{j} \in \mathcal{M}_{Y}$ for each $j$ and $S \subset \bigcup_{j=1}^{\infty}\left(A_{j} \times B_{j}\right)$.
6.50. Theorem. The set function $\varphi$ is an outer measure on $X \times Y$.

Proof. It is immediate from the definition that $\varphi \geq 0$ and $\varphi(\emptyset)=0$. To see that $\varphi$ is countably subadditive suppose $S \subset \bigcup_{k=1}^{\infty} S_{k}$ and assume that $\varphi\left(S_{k}\right)<\infty$ for each $k$. Fix $\varepsilon>0$. Then for each $k$ there is a sequence $\left\{A_{j}^{k} \times B_{j}^{k}\right\}_{j=1}^{\infty}$ with $A_{j}^{k} \in \mathcal{M}_{X}$ and $B_{j}^{k} \in \mathcal{M}_{Y}$ for each $j$ such that

$$
S_{k} \subset \bigcup_{j=1}^{\infty}\left(A_{j}^{k} \times B_{j}^{k}\right)
$$

and

$$
\sum_{j=1}^{\infty} \mu\left(A_{j}^{k}\right) \nu\left(B_{j}^{k}\right)<\varphi\left(S_{k}\right)+\frac{\varepsilon}{2^{k}}
$$

Thus

$$
\begin{aligned}
\varphi(S) & \leq \sum_{k=1}^{\infty} \sum_{j=1}^{\infty} \mu\left(A_{j}^{k}\right) \nu\left(B_{j}^{k}\right) \\
& \leq \sum_{k=1}^{\infty}\left(\varphi\left(S_{k}\right)+\frac{\varepsilon}{2^{k}}\right) \\
& \leq \sum_{k=1}^{\infty} \varphi\left(S_{k}\right)+\varepsilon
\end{aligned}
$$

for any $\varepsilon>0$.

Since $\varphi$ is an outer measure, we know its measurable sets form a $\sigma$-algebra (See Corollary 4.11) which we denote by $\mathcal{M}_{X \times Y}$. Also, we denote by $\mu \times \nu$ the restriction of $\varphi$ to $\mathcal{M}_{X \times Y}$. The main objective of this section is to show that $\mu \times \nu$ may appropriately be called the "product measure" corresponding to $\mu$ and $\nu$, and that the integral of a function over $X \times Y$ with respect to $\mu \times \nu$ can be computed by iterated integration. This is the thrust of the next result.
6.51. Theorem (Fubini's Theorem). Suppose $\left(X, \mathcal{M}_{X}, \mu\right)$ and $\left(Y, \mathcal{M}_{Y}, \nu\right)$ are complete measure spaces.
(i) If $A \in \mathcal{M}_{X}$ and $B \in \mathcal{M}_{Y}$, then $A \times B \in \mathcal{M}_{X \times Y}$ and

$$
(\mu \times \nu)(A \times B)=\mu(A) \nu(B)
$$

(ii) If $S \in \mathcal{M}_{X \times Y}$ and $S$ is $\sigma$-finite with respect to $\mu \times \nu$ then

$$
\begin{aligned}
S_{y} & =\{x:(x, y) \in S\} \in \mathcal{M}_{X}, \quad \text { for } \quad \nu-\text { a.e. } y \in Y \\
S_{x} & =\{y:(x, y) \in S\} \in \mathcal{M}_{Y}, \quad \text { for } \quad \mu-\text { a.e. } x \in X, \\
y & \mapsto \mu\left(S_{y}\right) \quad \text { is } \mathcal{M}_{Y} \text {-measurable, } \\
x & \mapsto \nu\left(S_{x}\right) \quad \text { is } \mathcal{M}_{X} \text {-measurable, } \\
(\mu \times \nu)(S) & =\int_{X} \nu\left(S_{x}\right) d \mu(x)=\int_{X}\left[\int_{Y} \chi_{S}(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{Y} \mu\left(S_{y}\right) d \nu(y)=\int_{Y}\left[\int_{X} \chi_{S}(x, y) d \mu(x)\right] d \nu(y) .
\end{aligned}
$$

(iii) If $f \in L^{1}\left(X \times Y, \mathcal{M}_{X \times Y}, \mu \times \nu\right)$, then

$$
\begin{aligned}
y & \mapsto f(x, y) \quad \text { is } \nu \text {-integrable for } \mu \text {-a.e. } x \in X, \\
x & \mapsto f(x, y) \quad \text { is } \mu \text {-integrable for } \nu \text {-a.e. } y \in Y, \\
x & \mapsto \int_{Y} f(x, y) d \nu(y) \quad \text { is } \mu \text {-integrable, } \\
y & \mapsto \int_{X} f(x, y) d \mu(x) \quad \text { is } \nu \text {-integrable, } \\
\int_{X \times Y} f d(\mu \times \nu) & =\int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{Y}\left[\int_{X} f(x, y) d \mu(x)\right] d \nu(y) .
\end{aligned}
$$

Proof. Let $\mathcal{F}$ denote the collection of all subsets $S$ of $X \times Y$ such that

$$
S_{x}:=\{y:(x, y) \in S\} \in \mathcal{M}_{Y} \quad \text { for } \mu \text {-a.e. } x \in X
$$

and that the function

$$
x \mapsto \nu\left(S_{x}\right) \quad \text { is } \mathcal{M}_{X} \text {-measurable. }
$$

For $S \in \mathcal{F}$ set

$$
\rho(S)=\int_{X} \nu\left(S_{x}\right) d \mu(x)=\int_{X}\left[\int_{Y} \chi_{S}(x, y) d \nu(y)\right] d \mu(x) .
$$

Another words, $\mathcal{F}$ is precisely the family of sets that makes is possible to define $\rho$. Proof of (i). First, note that $\rho$ is monotone on $\mathcal{F}$. Next observe that if $\cup_{j=1}^{\infty} S_{j}$ is a countable union of disjoint $S_{j} \in \mathcal{F}$, then clearly $\cup_{j=1}^{\infty} S_{j} \in \mathcal{F}$ and the Monotone Convergence Theorem implies

$$
\begin{equation*}
\sum_{j=1}^{\infty} \rho\left(S_{j}\right)=\rho\left(\bigcup_{j=1}^{\infty} S_{j}\right) ; \quad \text { hence } \bigcup_{i=1}^{\infty} S_{i} \in \mathcal{F} \text {. } \tag{6.31}
\end{equation*}
$$

Finally, if $S_{1} \supset S_{2} \supset \cdots$ are members of $\mathcal{F}$ then

$$
\begin{equation*}
\bigcap_{j=1}^{\infty} S_{j} \in \mathcal{F} \tag{6.32}
\end{equation*}
$$

and if $\rho\left(S_{1}\right)<\infty$, then Lebesgue's Dominated Convergence Theorem yields

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \rho\left(S_{j}\right)=\rho\left(\bigcap_{j=1}^{\infty} S_{j}\right) ; \quad \text { hence } \bigcap_{j=1}^{\infty} S_{j} \in \mathcal{F} \text {. } \tag{6.33}
\end{equation*}
$$

Set

$$
\begin{aligned}
& P_{0}=\left\{A \times B: A \in \mathcal{M}_{X} \quad \text { and } \quad B \in \mathcal{M}_{Y}\right\} \\
& P_{1}=\left\{\bigcup_{j=1}^{\infty} S_{j}: S_{j} \in P_{0} \quad \text { for } \quad j=1,2, \ldots\right\} \\
& P_{2}=\left\{\bigcap_{j=1}^{\infty} S_{j}: S_{j} \in P_{1} \quad \text { for } \quad j=1,2, \ldots\right\}
\end{aligned}
$$

Note that if $A \in \mathcal{M}_{X}$ and $B \in \mathcal{M}_{Y}$, then $A \times B \in \mathcal{F}$ and

$$
\begin{equation*}
\rho(A \times B)=\mu(A) \nu(B) \tag{6.34}
\end{equation*}
$$

and thus $P_{0} \subset \mathcal{F}$. If $A_{1} \times B_{1}, A_{2} \times B_{2} \subset X \times Y$, then

$$
\begin{equation*}
\left(A_{1} \times B_{1}\right) \cap\left(A_{2} \times B_{2}\right)=\left(A_{1} \cap A_{2}\right) \times\left(B_{1} \cap B_{2}\right) \tag{6.35}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(A_{1} \times B_{1}\right) \backslash\left(A_{2} \times B_{2}\right)=\left(\left(A_{1} \backslash A_{2}\right) \times B_{1}\right) \cup\left(\left(A_{1} \cap A_{2}\right) \times\left(B_{1} \backslash B_{2}\right)\right) \tag{6.36}
\end{equation*}
$$

It follows from Lemma $4.7,(6.35)$ and (6.36) that each member of $P_{1}$ can be written as a countable disjoint union of members of $P_{0}$ and since $\mathcal{F}$ is closed under countable disjoint unions, we have $P_{1} \subset \mathcal{F}$. It also follows from (6.35) that any finite intersection of members of $P_{1}$ is also a member of $P_{1}$. Therefore, from (6.32), $P_{2} \subset \mathcal{F}$. In summary, we have

$$
\begin{equation*}
P_{0}, P_{1}, P_{2} \subset \mathcal{F} \tag{6.37}
\end{equation*}
$$

Suppose $S \subset X \times Y,\left\{A_{j}\right\} \subset \mathcal{M}_{X},\left\{B_{j}\right\} \subset \mathcal{M}_{Y}$ and $S \subset R=\cup_{j=1}^{\infty}\left(A_{j} \times B_{j}\right)$. Using (6.34) and that $R \in P_{1} \subset \mathcal{F}$, we obtain

$$
\rho(R) \leq \sum_{j=1}^{\infty} \rho\left(A_{j} \times B_{j}\right)=\sum_{j=1}^{\infty} \mu\left(A_{j}\right) \nu\left(B_{j}\right)
$$

Thus, by the definition of $\sigma,(6.30)$,

$$
\begin{equation*}
\inf \left\{\rho(R): S \subset R \in P_{1}\right\} \leq \sigma(S) \tag{6.38}
\end{equation*}
$$

To establish the opposite inequality, note that if $S \subset R=\cup_{j=1}^{\infty}\left(A_{j} \times B_{j}\right)$ where the sets $A_{j} \times B_{j}$ are disjoint, then referring to (6.31)

$$
\sigma(S) \leq \sum_{j=1}^{\infty} \mu\left(A_{j}\right) \nu\left(B_{j}\right)=\rho(R)
$$

and consequently, with (6.38), we have

$$
\begin{equation*}
\sigma(S)=\inf \left\{\rho(R): S \subset R \in P_{1}\right\} \tag{6.39}
\end{equation*}
$$

for each $S \subset X \times Y$. If $A \in \mathcal{M}_{X}$ and $B \in \mathcal{M}_{Y}$, then $A \times B \in P_{0} \subset \mathcal{F}$ and hence, for any $R \in P_{1}$ with $A \times B \subset R$.

$$
\begin{aligned}
\sigma(A \times B) & \leq \mu(A) \nu(B) & & \text { by }(6.30) \\
& =\rho(A \times B) & & \text { by }(6.34) \\
& \leq \rho(R) & & \text { because } \rho \text { is monotone. }
\end{aligned}
$$

Therefore, by (6.39) and (6.34)

$$
\begin{equation*}
\sigma(A \times B)=\rho(A \times B)=\mu(A) \nu(B) \tag{6.40}
\end{equation*}
$$

Moreover if $T \subset R \in P_{1} \subset \mathcal{F}$, then using the additivity of $\rho$ (see(6.31)) it follows that

$$
\begin{array}{ll}
\sigma(T \backslash(A \times B))+\sigma(T \cap(A \times B)) & \\
\leq \rho(R \backslash(A \times B))+\rho(R \cap(A \times B)) & \text { by }(6.39) \\
=\rho(R) & \text { since } \rho \text { is additive. }
\end{array}
$$

In view of (6.39) we see that $\sigma(T \backslash(A \times B))+\sigma(T \cap(A \times B))$ for any $T \subset X \times Y$ and thus, (see Definition 4.3) that $A \times B$ is $\sigma$-measurable; that is, $A \times B \in \mathcal{M}_{X \times Y}$. Thus assertion (i) is proved.

Proof of (ii). Suppose $S \subset X \times Y$ and $\sigma(S)<\infty$. Then there is a sequence $\left\{R_{j}\right\} \subset P_{1}$ such that $S \subset R_{j}$ for each $j$ and

$$
\begin{equation*}
\sigma(S)=\lim _{j \rightarrow \infty} \rho\left(R_{j}\right) \tag{6.41}
\end{equation*}
$$

Set

$$
R=\bigcap_{j=1}^{\infty} R_{j} \in P_{2}
$$

Since $P_{2} \subset \mathcal{F}$ and $\sigma(S)<\infty$ the Dominated Convergence Theorem implies

$$
\begin{equation*}
\rho(R)=\lim _{m \rightarrow \infty} \rho\left(\bigcap_{j=1}^{m} R_{j}\right) . \tag{6.42}
\end{equation*}
$$

Thus, since $S \subset R \subset \cap_{j=1}^{m} R_{j} \in P_{2}$ for each finite $m$, by (6.42) we have

$$
\sigma(S) \leq \lim _{m \rightarrow \infty} \rho\left(\bigcap_{j=1}^{m} R_{j}\right)=\rho(R) \leq \lim _{m \rightarrow \infty} \rho\left(R_{m}\right)=\sigma(S)
$$

which implies that
(6.43) for each $S \subset X \times Y$ there is $R \in P_{2}$ such that $S \subset R$ and $\sigma(S)=\rho(R)$.

We now are in a position to finish the proof of assertion (ii). First suppose $S \subset X \times Y, S \subset R \in P_{2}$ and $\rho(R)=0$. Then $\nu\left(R_{x}\right)=0$ for $\mu$-a.e. $x \in X$ and
$S_{x} \subset R_{x}$ for each $x \in X$. Since $\nu$ is complete, $S_{x} \in \mathcal{M}_{Y}$ for $\mu$-a.e. $x \in X$ and $S \in \mathcal{F}$ with $\rho(S)=0$. In particular we see that if $S \subset X \times Y$ with $\sigma(S)=0$, then $S \in \mathcal{F}$ and $\rho(S)=0$.

Now suppose $S \in \mathcal{M}_{X \times Y}$ and $(\mu \times \nu)(S)<\infty$. Then, from (6.43), there is an $R \in P_{2}$ such that $S \subset R$ and

$$
(\mu \times \nu)(S)=\sigma(S)=\rho(R)
$$

From assertion (i) we see that $R \in \mathcal{M}_{X \times Y}$, and since $(\mu \times \nu)(S)<\infty$

$$
(\mu \times \nu)(R \backslash S)=0
$$

This in turn implies that $R \backslash S \in \mathcal{F}$ and $\rho(R \backslash S)=0$. Since $\nu$ is complete and

$$
R_{x} \backslash S_{x} \in \mathcal{M}_{Y}
$$

for $\mu$-a.e. $x \in X$ we see that $S_{x} \in \mathcal{M}_{Y}$ for $\mu$-a.e. $x \in X$ and thus that $S \in \mathcal{F}$ with

$$
(\mu \times \nu)(S)=\rho(S)=\int_{X} \nu\left(S_{x}\right) d \mu(x)
$$

If $S \in \mathcal{M}_{X \times Y}$ is $\sigma$-finite with respect the measure $\mu \times \nu$, then there exists a sequence $\left\{S_{j}\right\}$ of disjoint sets $S_{j} \in \mathcal{M}_{X \times Y}$ with $(\mu \times \nu)\left(S_{j}\right)<\infty$ for each $j$ such that

$$
S=\bigcup_{j=1}^{\infty} S_{j} .
$$

Since the sets are disjoint and each $S_{j} \in \mathcal{F}$ we have $S \in \mathcal{F}$ and

$$
(\mu \times \nu)(S)=\sum_{j=1}^{\infty}(\mu \times \nu)\left(S_{j}\right)=\sum_{j=1}^{\infty} \rho\left(S_{j}\right)=\rho(S) .
$$

Of course the above argument remains valid if the roles of $\mu$ and $\nu$ are interchanged, and thus we have proved assertion (ii).

Proof of (iii). Assume first that $f \in L^{1}\left(X \times Y, \mathcal{M}_{X \times Y}, \mu \times \nu\right)$ and $f \geq 0$. Fix $t>1$ and set

$$
E_{k}=\left\{(x, y): t^{k}<f(x, y) \leq t^{k+1}\right\} .
$$

for each $k=0, \pm 1, \pm 2, \ldots$. Then each $E_{k} \in \mathcal{M}_{X \times Y}$ with $(\mu \times \nu)\left(E_{k}\right)<\infty$. In view of (ii) the function

$$
f_{t}=\sum_{k=-\infty}^{\infty} t^{k} \chi_{E_{k}}
$$

satisfies the first four assertions of (ii) and

$$
\begin{align*}
\int f_{t} d(\mu \times \nu) & =\sum_{k=-\infty}^{\infty} t^{k}(\mu \times \nu)\left(E_{k}\right) \\
& =\sum_{k=-\infty}^{\infty} t^{k} \rho\left(E_{k}\right)  \tag{6.40}\\
& =\sum_{k=-\infty}^{\infty} t^{k} \int_{X}\left[\int_{Y} \chi_{E_{k}}(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{X}\left[\sum_{k=-\infty}^{\infty} t^{k} \int_{Y} \chi_{E_{k}}(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{X}\left[\int_{Y} f_{t}(x, y) d \nu(y)\right] d \mu(x)
\end{align*}
$$

by the Monotone Convergence Theorem. Similarly

$$
\int f_{t} d(\mu \times \nu)=\int_{Y}\left[\int_{X} f_{t}(x, y) d \mu(x)\right] d \nu(y)
$$

Since

$$
\frac{1}{t} f \leq f_{t} \leq f
$$

we see that $f_{t}(x, y) \rightarrow f(x, y)$ as $t \rightarrow 1^{+}$for each $(x, y) \in X \times Y$. Thus the function

$$
y \mapsto f(x, y)
$$

is $\mathcal{M}_{Y}$-measurable for $\mu$-a.e. $x \in X$. It follows that

$$
\frac{1}{t} \int_{Y} f(x, y) d \nu(y) \leq \int_{Y} f_{t}(x, y) d \nu(y) \leq \int_{Y} f(x, y) d \nu(y)
$$

for $\mu$-a.e. $x \in X$, the function

$$
x \mapsto \int_{Y} f(x, y) d \nu(y)
$$

is $\mathcal{M}_{X}$-measurable and

$$
\begin{aligned}
\frac{1}{t} \int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x) & \leq \int_{X}\left[\int_{Y} f_{t}(x, y) d \nu(y)\right] d \mu(x) \\
& \leq \int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x)
\end{aligned}
$$

Thus we see

$$
\begin{aligned}
\int f d(\mu \times \nu) & =\lim _{t \rightarrow 1^{+}} \int f_{t}(x, y) d(\mu \times \nu) \\
& =\lim _{t \rightarrow 1^{+}} \int_{X}\left[\int_{Y} f_{t}(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x)
\end{aligned}
$$

Since the first integral above is finite we have established the first and third parts of assertion (iii) as well as the first half of the fifth part. The remainder of (iii) follows by an analogous argument.

To extend the proof to general $f \in L^{1}\left(X \times Y, \mathcal{M}_{X \times Y}, \mu \times \nu\right)$ we need only recall that

$$
f=f^{+}-f^{-}
$$

where $f^{+}$and $f^{-}$are nonnegative integrable functions.
It is important to observe that the hypothesis (iii) in Fubini's theorem, namely that $f \in L^{1}\left(X \times Y, \mathcal{M}_{X \times Y}, \mu \times \nu\right)$, is necessary. Indeed, consider the following example.
6.52. Example. Let $Q$ denote the unit square $[0,1] \times[0,1]$ and consider a sequence of subsquares $Q_{k}$ defined as follows: Let $Q_{1}:=[0,1 / 2] \times[0,1 / 2]$. Let $Q_{2}$ be a square with half the area of $Q_{1}$ and placed so that $Q_{1} \cap Q_{2}=\{(1 / 2,1 / 2)\}$; that is, so that its "southwest" vertex is the same as the "northeast" vertex of $Q_{1}$. Similarly, let $Q_{3}$ be a square with half the area of $Q_{2}$ and as before, place it so that its "southwest" vertex is the same as the "northeast" vertex of $Q_{2}$. In this way, we obtain a sequence of squares $\left\{Q_{k}\right\}$ all of whose southwest-northeast diagonal vertices lie on the line $y=x$. Subdivide each subsquare $Q_{k}$ into four equal squares $Q_{k}^{(1)}, Q_{k}^{(2)}, Q_{k}^{(3)}, Q_{k}^{(4)}$, where we will regard $Q_{k}^{(1)}$ as occupying the "first quadrant", $Q_{k}^{(2)}$ the "second quadrant", $Q_{k}^{(3)}$ the "third quadrant" and $Q_{k}^{(4)}$ the "fourth quadrant." In the next section it will be shown that two dimensional Lebesgue measure $\lambda_{2}$ is the same as the product $\lambda_{1} \times \lambda_{1}$ where $\lambda_{1}$ denotes onedimensional Lebesgue measure. Define a function $f$ on $Q$ such that $f=0$ on complement of the $Q_{k}$ 's and otherwise, on each $Q_{k}$ define $f=\frac{1}{\lambda_{2}\left(Q_{k}\right)}$ on subsquares in the first and third quadrants and $f=-\frac{1}{\lambda_{2}\left(Q_{k}\right)}$ on the subsquares in the second and fourth quadrants. Clearly,

$$
\int_{Q_{k}}|f| d \lambda_{2}=1
$$

and therefore

$$
\int_{Q}|f| d \lambda_{2}=\sum_{k} \int_{Q_{k}}|f|=\infty
$$

whereas

$$
\int_{0}^{1} f(x, y) d \lambda_{1}(x)=\int_{0}^{1} f(x, y) d \lambda_{1}(y)=0
$$

This is an example where the iterated integral exists but Fubini's Theorem does not hold because $f$ is not integrable. However, this integrability hypothesis is not
necessary if $f \geq 0$ and if $f$ is measurable in each variable separately. The proof of this follows readily from the proof of Theorem 6.51.
6.53. Corollary (Tonelli). If $f$ is a nonnegative $\mathcal{M}_{X \times Y}$-measurable function and $\{(x, y): f(x, y) \neq 0\}$ is $\sigma$-finite with respect to the measure $\mu \times \nu$, then the function

$$
\begin{aligned}
y \mapsto f(x, y) & \text { is } \mathcal{M}_{Y} \text {-measurable for } \mu \text {-a.e. } x \in X \\
x \mapsto f(x, y) & \text { is } \mathcal{M}_{X} \text {-measurable for } \nu \text {-a.e. } y \in Y, \\
x & \mapsto \int_{Y} f(x, y) d \nu(y) \quad \text { is } \mathcal{M}_{X} \text {-measurable } \\
y & \mapsto \int_{X} f(x, y) d \mu(x) \quad \text { is } \mathcal{M}_{Y} \text {-measurable }
\end{aligned}
$$

and

$$
\int_{X \times Y} f d(\mu \times \nu)=\int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x)=\int_{Y}\left[\int_{X} f(x, y) d \mu(x)\right] d \nu(y)
$$

in the sense that either both expressions are infinite or both are finite and equal.

Proof. Let $\left\{f_{k}\right\}$ be a sequence of nonnegative real-valued measurable functions with finite range such that $f_{k} \leq f_{k+1}$ and $\lim _{k \rightarrow \infty} f_{k}=f$. By assertion (ii) of Theorem 6.51 the conclusion of the corollary holds for each $f_{k}$. For each $k$ let $N_{k}$ be a $\mathcal{M}_{X}$-measurable subset of $X$ such that $\mu\left(N_{k}\right)=0$ and

$$
y \mapsto f_{k}(x, y)
$$

is $\mathcal{M}_{Y}$-measurable for each $x \in X-N_{k}$. Set $N=\cup_{k=1}^{\infty} N_{k}$. Then $\mu(N)=0$ and for each $x \in X-N$

$$
y \mapsto f(x, y)=\lim _{k \rightarrow \infty} f_{k}(x, y)
$$

is $\mathcal{M}_{Y}$-measurable and by the Monotone Convergence Theorem

$$
\begin{equation*}
\int_{Y} f(x, y) d \nu(y)=\lim _{k \rightarrow \infty} \int_{Y} f_{k}(x, y) d \nu(y) \tag{6.44}
\end{equation*}
$$

for $x \in X-N$.

Theorem 6.51 implies that $h_{k}(x):=\int_{Y} f_{k}(x, y) d \nu(y)$ is $\mathcal{M}_{X}$-measurable. Since $0 \leq h_{k} \leq h_{k+1}$, we can use again the Monotone Convergence Theorem to obtain

$$
\begin{aligned}
\int_{X \times Y} f d(\mu \times \nu) & =\lim _{k \rightarrow \infty} \int_{X \times Y} f_{k}(x, y) d(\mu \times \nu) \\
& =\lim _{k \rightarrow \infty} \int_{X}\left[\int_{Y} f_{k}(x, y) d \nu(y)\right] d \mu(x) \\
& =\lim _{k \rightarrow \infty} \int_{X} h_{k}(x) d \mu(x) \\
& =\int_{X} \lim _{k \rightarrow \infty} h_{k}(x) d \mu(x) \\
& =\int_{X} \lim _{k \rightarrow \infty}\left[\int_{Y} f_{k}(x, y) d \nu(y)\right] d \mu(x) \\
& =\int_{X}\left[\int_{Y} f(x, y) d \nu(y)\right] d \mu(x),
\end{aligned}
$$

where the last line follows from (6.44). The reversed iteration can be obtained with a similar argument.

### 6.10. Lebesgue Measure as a Product Measure

We will now show that $n$-dimensional Lebesgue measure on $\mathbb{R}^{n}$ is a product of lower dimensional Lebesgue measures.

For each positive integer $k$ let $\lambda_{k}$ denote Lebesgue measure on $\mathbb{R}^{k}$ and let $\mathcal{M}_{k}$ denote the $\sigma$-algebra of Lebesgue measurable subsets of $\mathbb{R}^{k}$.
6.54. Theorem. For each pair of positive integers $n$ and $m$

$$
\lambda_{n+m}=\lambda_{n} \times \lambda_{m}
$$

Proof. Let $\varphi$ denote the outer measure on $\mathbb{R}^{n} \times \mathbb{R}^{m}$ defined as in Definition 6.49 with $\mu=\lambda_{n}$ and $\nu=\lambda_{m}$. We will show that $\varphi=\lambda_{n+m}^{*}$.

If $A \in \mathcal{M}_{n}$ and $B \in \mathcal{M}_{m}$ are bounded sets and $\varepsilon>0$, then there are open sets $U \supset A, V \supset B$ such that

$$
\begin{gathered}
\lambda_{n}(U-A)<\varepsilon \\
\lambda_{m}(V-B)<\varepsilon
\end{gathered}
$$

and hence

$$
\begin{equation*}
\lambda_{n}(U) \lambda_{m}(V) \leq \lambda_{n}(A) \lambda_{m}(B)+\varepsilon\left(\lambda_{n}(A)+\lambda_{m}(B)\right)+\varepsilon^{2} \tag{6.45}
\end{equation*}
$$

Suppose $E$ is a bounded subset of $\mathbb{R}^{n+m}$ and $\left\{A_{k} \times B_{k}\right\}_{k=1}^{\infty}$ a sequence of subsets of $\mathbb{R}^{n} \times \mathbb{R}^{m}$ such that $A_{k} \in \mathcal{M}_{n}, B_{k} \in \mathcal{M}_{m}$, and $E \subset \cup_{k=1}^{\infty} A_{k} \times B_{k}$. Assume that the sequences $\left\{\lambda_{n}\left(A_{k}\right)\right\}_{k=1}^{\infty}$ and $\left\{\lambda_{m}\left(B_{k}\right)\right\}_{k=1}^{\infty}$ are bounded. Fix $\varepsilon>0$. In view of (6.45) there exist open sets $U_{k}, V_{k}$ such that $A_{k} \subset U_{k}, B_{k} \subset V_{k}$, and

$$
\sum_{k=1}^{\infty} \lambda_{n}\left(A_{k}\right) \lambda_{m}\left(B_{k}\right) \geq \sum_{k=1}^{\infty} \lambda_{n}\left(U_{k}\right) \lambda_{m}\left(V_{k}\right)-\varepsilon
$$

It is not difficult to show that each of the open sets $U_{k} \times V_{k}$ can be written as a countable union of nonoverlapping closed intervals

$$
U_{k} \times V_{k}=\bigcup_{l=1}^{\infty} I_{l}^{k} \times J_{l}^{k}
$$

where $I_{l}^{k}, J_{l}^{k}$ are closed intervals in $\mathbb{R}^{n}, \mathbb{R}^{m}$ respectively. Thus for each $k$

$$
\lambda_{n}\left(U_{k}\right) \lambda_{m}\left(V_{k}\right)=\sum_{l=1}^{\infty} \lambda_{n}\left(I_{l}^{k}\right) \lambda_{m}\left(J_{l}^{k}\right)=\sum_{l=1}^{\infty} \lambda_{n+m}\left(I_{l}^{k} \times J_{l}^{k}\right)
$$

It follows that

$$
\sum_{k=1}^{\infty} \lambda_{n}\left(A_{k}\right) \lambda_{m}\left(B_{k}\right) \geq \sum_{k=1}^{\infty} \lambda_{n}\left(U_{k}\right) \lambda_{m}\left(V_{k}\right)-\varepsilon \geq \lambda_{n+m}^{*}(E)-\varepsilon
$$

and hence

$$
\varphi(E) \geq \lambda_{n+m}^{*}(E)
$$

whenever $E$ is a bounded subset of $\mathbb{R}^{n+m}$.
In case $E$ is an unbounded subset of $\mathbb{R}^{n+m}$ we have

$$
\varphi(E) \geq \lambda_{n+m}^{*}(E \cap B(0, j))
$$

for each positive integer $j$. Since $\lambda_{n+m}^{*}$ is a Borel regular outer measure (see Exercise 9, Section 4.3), there is a Borel set $A_{j} \supset E \cap B(0, j)$ such that $\lambda\left(A_{j}\right)=\lambda_{m+n}^{*}(E \cap$ $B(0, j))$. With $A:=\cup_{j=1}^{\infty} A_{j}$, we have

$$
\lambda_{n+m}^{*}(E) \leq \lambda_{n+m}(A)=\lim _{j \rightarrow \infty} \lambda_{n+m}\left(A_{j}\right)=\lim _{j \rightarrow \infty} \lambda_{n+m}^{*}(E \cap B(0, j)) \leq \lambda_{n+m}^{*}(E)
$$

and therefore

$$
\lim _{j \rightarrow \infty} \lambda_{m+n}^{*}(E \cap B(0, j))=\lambda_{m+n}^{*}(E)
$$

This yields

$$
\varphi(E) \geq \lambda_{m+n}^{*}(E)
$$

for each $E \subset \mathbb{R}^{n+m}$.
On the other hand it is immediate from the definitions of the two outer measures that

$$
\varphi(E) \leq \lambda_{m+n}^{*}(E)
$$

for each $E \subset \mathbb{R}^{n+m}$.

### 6.11. Convolution

As an application of Fubini's Theorem, we determine conditions on functions $f$ and $g$ that ensure the existence of the convolution $f * g$ and deduce the basic properties of convolution.
6.55. Definition. Given two Lebesgue measurable functions $f$ and $g$ on $\mathbb{R}^{n}$ we define the convolution $f * g$ of $f$ and $g$ to be the function defined for each $x \in \mathbb{R}^{n}$ by

$$
(f * g)(x)=\int_{\mathbb{R}^{n}} f(y) g(x-y) d y
$$

Here and in the remainder of this section we will indicate integration with respect to Lebesgue measure by $d x, d y$, etc.

We first observe that if $g$ is a nonnegative Lebesgue measurable function on $\mathbb{R}^{n}$, then

$$
\int_{\mathbb{R}^{n}} g(x-y) d y=\int_{\mathbb{R}^{n}} g(y) d y
$$

for any $x \in \mathbb{R}^{n}$. This follows readily from the definition of the integral and the fact that $\lambda_{n}$ is invariant under translation.

To study the integrability properties of the convolution of two functions we will need the following lemma.
6.56. Lemma. If $f$ is a Lebesgue measurable function on $\mathbb{R}^{n}$, then the function $F$ defined on $\mathbb{R}^{n} \times \mathbb{R}^{n}=\mathbb{R}^{2 n}$ by

$$
F(x, y)=f(x-y)
$$

is $\lambda_{2 n}$-measurable.

Proof. First, define $F_{1}: \mathbb{R}^{2 n} \rightarrow \mathbb{R}$ by $F_{1}(x, y):=f(x)$ and observe that $F_{1}$ is $\lambda_{2 n}$-measurable because for any Borel set $B \subset \mathbb{R}$, we have $F_{1}^{-1}(B)=f^{-1}(B) \times$ $\mathbb{R}^{n}$. Then define $T: \mathbb{R}^{2 n} \rightarrow \mathbb{R}^{2 n}$ by $T(x, y)=(x-y, x+y)$ and note that $T^{-1}(x, y)=\left(\frac{x+y}{2}, \frac{y-x}{2}\right)$. The mean value theorem implies $\left|T\left(x_{1}, y_{1}\right)-T\left(x_{2}, y_{2}\right)\right| \leq$ $n^{2}\left|\left(x_{1}, y_{1}\right)-\left(x_{2}, y_{2}\right)\right|$ and $\left|T^{-1}\left(x_{1}, y_{1}\right)-T^{-1}\left(x_{2}, y_{2}\right)\right| \leq \frac{1}{2} n^{2}\left|\left(x_{1}, y_{1}\right)-\left(x_{2}, y_{2}\right)\right|$ for every $\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right) \in \mathbb{R}^{2 n}$. Therefore $T$ and $T^{-1}$ are Lipschitz functions in $\mathbb{R}^{2 n}$. Hence, it follows that $F_{1} \circ T=F$ is $\lambda_{2 n}$-measurable. Indeed, if $B \subset \mathbb{R}$ is a Borel set, then $E:=F_{1}^{-1}(B)$ is $\lambda_{2 n}$-measurable and thus can be expressed as $E=B_{1} \cup N$ where $B_{1} \subset \mathbb{R}^{2 n}$ is a Borel set and $\lambda_{2 n}(N)=0$. Consequently, $T^{-1}(E)=T^{-1}\left(B_{1}\right) \cup T^{-1}(N)$, which is the union of a Borel set and a set of $\lambda_{2 n^{-}}$ measure zero (see exercise 4.11).

We now prove a basic result concerning convolutions. Recall our notation $L^{p}\left(\mathbb{R}^{n}\right)$ for $L^{p}\left(\mathbb{R}^{n}, \mathcal{M}_{n}, \lambda_{n}\right)$ and $\|f\|_{p}$ for $\|f\|_{p, \mathbb{R}^{n} ; \lambda_{n}}$.
6.57. Theorem. If $f \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p \leq \infty$ and $g \in L^{1}\left(\mathbb{R}^{n}\right)$, then $f * g \in$ $L^{p}\left(\mathbb{R}^{n}\right)$ and

$$
\|f * g\|_{p} \leq\|f\|_{p}\|g\|_{1}
$$

Proof. Observe that $|f * g| \leq|f| *|g|$ and thus it suffices to prove the assertion for $f, g \geq 0$. Then by Lemma 6.56 the function

$$
(x, y) \mapsto f(y) g(x-y)
$$

is nonnegative and $\mathcal{M}_{2 n}$-measurable and by Corollary 6.53

$$
\begin{aligned}
\int(f * g)(x) d x & =\iint f(y) g(x-y) d y d x \\
& =\int f(y)\left[\int g(x-y) d x\right] d y \\
& =\int f(y) d y \int g(x) d x
\end{aligned}
$$

Thus the assertion holds if $p=1$.
In case $p=\infty$ we see

$$
(f * g)(x) \leq\|f\|_{\infty} \int g(x-y) d y=\|f\|_{\infty}\|g\|_{1}
$$

whence

$$
\|f * g\|_{\infty} \leq\|f\|_{\infty}\|g\|_{1}
$$

Finally suppose $1<p<\infty$. Then

$$
\begin{aligned}
(f * g)(x) & =\int f(y)(g(x-y))^{\frac{1}{p}}(g(x-y))^{1-\frac{1}{p}} d y \\
& \leq\left(\int f^{p}(y) g(x-y) d y\right)^{\frac{1}{p}}\left(\int g(x-y) d y\right)^{1-\frac{1}{p}} \\
& =\left(f^{p} * g\right)^{\frac{1}{p}}(x)\|g\|_{1}^{1-\frac{1}{p}}
\end{aligned}
$$

Thus

$$
\begin{aligned}
\int(f * g)^{p}(x) d x & \leq \int\left(f^{p} * g\right)(x) d x\|g\|_{1}^{p-1} \\
& =\left\|f^{p}\right\|_{1}\|g\|_{1}\|g\|_{1}^{p-1} \\
& =\|f\|_{p}^{p}\|g\|_{1}^{p}
\end{aligned}
$$

and the assertion is proved.

If we fix $g \in L^{1}\left(\mathbb{R}^{n}\right)$ and set

$$
T(f)=f * g,
$$

then we may interpret the theorem as saying that for any $1 \leq p \leq \infty$,

$$
T: L^{p}\left(\mathbb{R}^{n}\right) \rightarrow L^{p}\left(\mathbb{R}^{n}\right)
$$

is a bounded linear mapping. Such mappings induced by convolution will be further studied in Chapter 9.

## Exercises for Section 6.11

1. (a) For $p>1$ and $p^{\prime}:=p /(p-1)$, prove that if $f \in L^{p}\left(\mathbb{R}^{n}\right)$ and $g \in L^{p^{\prime}}\left(\mathbb{R}^{n}\right)$, then

$$
f * g(x) \leq\|f\|_{p}\|g\|_{p^{\prime}} .
$$

for any $x \in \mathbb{R}^{n}$.
(b) Suppose that $f \in L^{p}\left(\mathbb{R}^{n}\right)$ and $g \in L^{p^{\prime}}\left(\mathbb{R}^{n}\right)$. Prove that $f * g$ vanishes at infinity. That is, prove that for each $\varepsilon>0$, there exists $R>0$ such that

$$
f * g(x)<\varepsilon \text { for all }|x|>R .
$$

2. Let $\phi$ be a non-negative, real-valued function in $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ with the property that

$$
\int_{\mathbb{R}^{n}} \phi(x) d x=1, \operatorname{spt} \phi \subset \bar{B}(0,1) .
$$

An example of such a function is given by

$$
\phi(x)= \begin{cases}C \exp \left[-1 /\left(1-|x|^{2}\right)\right] & \text { if }|x|<1 \\ 0 & \text { if }|x| \geq 1\end{cases}
$$

where $C$ is chosen so that $\int_{R^{n}} \phi=1$. For $\varepsilon>0$, the function $\phi_{\varepsilon}(x):=\varepsilon^{-n} \phi(x / \varepsilon)$ belongs to $C_{0}^{\infty}\left(R^{n}\right)$ and $\operatorname{spt} \phi_{\varepsilon} \subset \bar{B}(0, \varepsilon)$. $\phi_{\varepsilon}$ is called a regularizer (or mollifier) and the convolution

$$
u_{\varepsilon}(x):=\phi_{\varepsilon} * u(x):=\int_{\mathbb{R}^{n}} \phi_{\varepsilon}(x-y) u(y) d y
$$

defined for functions $u \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ is called the regularization (mollification) of $u$. As a consequence of Fubini's theorem, we have

$$
\|u * v\|_{p} \leq\|u\|_{p}\|v\|_{1}
$$

whenever $1 \leq p \leq \infty, u \in L^{p}\left(\mathbb{R}^{n}\right)$ and $v \in L^{1}\left(\mathbb{R}^{n}\right)$.
Prove the following (see Theorem 10.1):
(a) If $u \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$, then for every $\varepsilon>0, u_{\varepsilon} \in C^{\infty}\left(\mathbb{R}^{n}\right)$.
(b) If $u$ is continuous, then $u_{\varepsilon}$ converges to $u$ uniformly on compact subsets of $\mathbb{R}^{n}$.
3. If $u \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<\infty$, then $u_{\varepsilon} \in L^{p}\left(\mathbb{R}^{n}\right),\left\|u_{\varepsilon}\right\|_{p} \leq\|u\|_{p}$, and $\lim _{\varepsilon \rightarrow 0}\left\|u_{\varepsilon}-u\right\|_{p}=0$.
4. Let $\mu$ be a Radon measure on $\mathbb{R}^{n}, x \in \mathbb{R}^{n}$ and $0<\alpha<n$. Then

$$
\int_{\mathbb{R}^{n}} \frac{d \mu(y)}{|x-y|^{n-\alpha}}=(n-\alpha) \int_{0}^{\infty} r^{\alpha-n-1} \mu(B(x, r)) d r
$$

provided that

$$
\int_{\mathbb{R}^{n}} \frac{d \mu(y)}{|x-y|^{n-\alpha}}<\infty
$$

5. In this problem, we will consider $\mathbb{R}^{2}$ for simplicity, but everything carries over to $\mathbb{R}^{n}$. Let $P$ be a polynomial in $\mathbb{R}^{2}$; that is, $P$ has the form
$P(x, y)=a_{n} x^{n} y^{n}+a_{n-1} x^{n} y^{n-1}+b_{n-1} x^{n-1} y^{n}+\cdots+a_{1} x^{1} y^{0}+b_{1} x y^{1}+a_{0}$,
where the $a^{\prime} s$ and $b^{\prime} s$ are real numbers and $n \in \mathbb{N}$. Let $\varphi_{\varepsilon}$ denote the mollifying kernel discussed in the previous problem. Prove that $\varphi_{\varepsilon} * P$ is also a polynomial. In other words, it isn't possible to make a polynomial anymore smooth than itself.
6. Consider $(X, \mathcal{M}, \mu)$ where $\mu$ is $\sigma$-finite and complete and suppose $f \in L^{1}(X)$ is nonnegative. Let

$$
G_{f}:=\{(x, y) \in X \times[0, \infty]: 0 \leq y \leq f(x)\}
$$

Prove the following:
(a) The set $G_{f}$ is $\mu \times \lambda_{1}$-measurable.
(b) $\mu \times \lambda_{1}\left(G_{f}\right)=\int_{X} f d \mu$.

This shows that the "area under the graph is the integral of the function."

### 6.12. Distribution Functions

Here we will study an interesting and useful connection between abstract integration and Lebesgue integration.

Let $(X, \mathcal{M}, \mu)$ be a complete $\sigma$-finite measure space. Let $f$ be a measurable function on $X$ and for $t \in \mathbb{R}$ set

$$
E_{t}=\{x:|f(x)|>t\} \in \mathcal{M}
$$

We have the following
6.58. Definition. The distribution function of $f$ is the nonincreasing function defined as

$$
A_{f}(t):=\mu\left(E_{t}\right)
$$

An interesting relation between $f$ and its distribution function can be deduced from Fubini's Theorem.
6.59. ThEOREM. If $f$ is nonnegative and measurable, then

$$
\begin{equation*}
\int_{X} f d \mu=\int_{[0, \infty)} A_{f} d \lambda=\int_{[0, \infty]} \mu(\{x: f(x)>t\}) d \lambda(t) \tag{6.46}
\end{equation*}
$$

Proof. Let $\widetilde{\mathcal{M}}$ denote the $\sigma$-algebra of measurable subsets of $X \times \mathbb{R}$ corresponding to $\mu \times \lambda$. Set

$$
W=\{(x, t): 0<t<f(x)\} \subset X \times \mathbb{R}
$$

Since $f$ is measurable there is a sequence $\left\{f_{k}\right\}$ of measurable simple functions such that $f_{k} \leq f_{k+1}$ and $\lim _{k \rightarrow \infty} f_{k}=f$ pointwise on $X$. If $f_{k}=\sum_{j=1}^{n_{k}} a_{j}^{k} \chi_{E_{j}^{k}}$ where for each $k$ the sets $\left\{E_{j}^{k}\right\}$ are disjoint and measurable, then

$$
W_{k}=\left\{(x, t): 0<t<f_{k}(x)\right\}=\bigcup_{j=1}^{n_{k}} E_{j}^{k} \times\left(0, a_{j}^{k}\right) \in \widetilde{\mathcal{M}}
$$

Since $\chi_{W}=\lim _{k \rightarrow \infty} \chi_{W_{k}}$ we see that $W \in \widetilde{\mathcal{M}}$. Thus by Corollary 6.53

$$
\begin{aligned}
\int_{[0, \infty)} A_{f} d \lambda & =\int_{\mathbb{R}} \int_{X} \chi_{W}(x, t) d \mu(x) d \lambda(t) \\
& =\int_{X} \int_{\mathbb{R}} \chi_{W}(x, t) d \lambda(t) d \mu(x) \\
& =\int_{X} \lambda(\{t: 0<t<f(x)\}) d \mu(x) \\
& =\int_{X} f d \mu
\end{aligned}
$$

Thus a nonnegative measurable function $f$ is integrable over $X$ with respect to $\mu$ if and only if its distribution function $A_{f}$ is integrable over $[0, \infty)$ with respect to one-dimensional Lebesgue measure $\lambda$.

If $\mu(X)<\infty$, then $A_{f}$ is a bounded monotone function and thus continuous $\lambda$ a.e. on $[0, \infty)$. In view of Theorem 6.19 this implies that $A_{f}$ is Riemann integrable on any compact interval in $[0, \infty)$ and thus that the right-hand side of (6.46) can be interpreted as an improper Riemann integral.

The simple idea behind the proof of Theorem 6.59 can readily be extended as in the following theorem.
6.60. Theorem. If $f$ is measurable and $1 \leq p<\infty$, then

$$
\int_{X}|f|^{p} d \mu=p \int_{[0, \infty)} t^{p-1} \mu(\{x:|f(x)|>t\}) d \lambda(t)
$$

Proof. Set

$$
W=\{(x, t): 0<t<|f(x)|\}
$$

and note that the function

$$
(x, t) \mapsto p t^{p-1} \chi_{W}(x, t)
$$

is $\widetilde{\mathcal{M}}$-measurable. Thus by Corollary 6.53

$$
\begin{aligned}
\int_{X}|f|^{p} d \mu & =\int_{X} \int_{(0,|f(x)|)} p t^{p-1} d \lambda(t) d \mu(x) \\
& =\int_{X} \int_{\mathbb{R}} p t^{p-1} \chi_{W}(x, t) d \lambda(t) d \mu(x) \\
& =\int_{\mathbb{R}} \int_{X} p t^{p-1} \chi_{W}(x, t) d \mu(x) d \lambda(t) \\
& =p \int_{[0, \infty)} t^{p-1} \mu(\{x:|f(x)|>t\}) d \lambda(t)
\end{aligned}
$$

6.61. Remark. A useful mnemonic relating to the previous result is that if $f$ is measurable and $1 \leq p<\infty$, then

$$
\int_{X}|f|^{p} d \mu=\int_{0}^{\infty} \mu(\{|f|>t\}) d t^{p} .
$$

## Exercises for Section 6.12

1. Suppose $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and let $A_{t}:=\{x:|f(x)|>t\}$. Prove that

$$
\lim _{t \rightarrow \infty} \int_{A_{t}}|f| d \lambda=0
$$

### 6.13. The Marcinkiewicz Interpolation Theorem

In the previous section, we employed Fubini's theorem extensively to investigate the properties of the distribution function. We close this chapter by pursuing this topic further to establish the Marcinkiewicz Interpolation Theorem, which has important applications in diverse areas of analysis, such as Fourier analysis and nonlinear potential theory. Later, in Chapter 7, we will see a beautiful interaction between this result and the Hardy-Littlewood Maximal function, Definiton 7.8.

In preparation for the main theorem of this section, we will need two preliminary results. The first, due to Hardy, gives two inequalities that are related to Jensen's inequality Exercise 10 , Section 6.5. If $f$ is a non-negative measurable function defined on the positive real numbers, let

$$
\begin{aligned}
& F(x)=\frac{1}{x} \int_{0}^{x} f(t) d t, x>0 \\
& G(x)=\frac{1}{x} \int_{x}^{\infty} f(t) d t, x>0
\end{aligned}
$$

Jensen's inequality states that, for $p \geq 1,[F(x)]^{p} \leq\|f\|_{p ;(0, x)}$ for each $x>0$ and thus provides an estimate of $F(x)^{p}$; Hardy's inequality (6.47) below, gives an estimate of a weighted integral of $F^{p}$.
6.62. Lemma (Hardy's Inequalities). If $1 \leq p<\infty, r>0$ and $f$ is a nonnegative measurable function on $(0, \infty)$, then with $F$ and $G$ defined as above,

$$
\begin{align*}
& \int_{0}^{\infty}[F(x)]^{p} x^{p-r-1} d x \leq\left(\frac{p}{r}\right)^{p} \int_{0}^{\infty}[f(t)]^{p} t^{p-r-1} d t  \tag{6.47}\\
& \int_{0}^{\infty}[G(x)]^{p} x^{p+r-1} d x \leq\left(\frac{p}{r}\right)^{p} \int_{0}^{\infty}[f(t)]^{p} t^{p+r-1} d t \tag{6.48}
\end{align*}
$$

Proof. To prove (6.47), we apply Jensen's inequality (Exercise 10, Section 6.5 ) with the measure $t^{(r / p)-1} d t$, we obtain

$$
\begin{align*}
\left(\int_{0}^{x} f(t) d t\right)^{p} & =\left(\int_{0}^{x} f(t) t^{1-(r / p)} t^{(r / p)-1} d t\right)^{p}  \tag{6.49}\\
& \leq\left(\frac{p}{r}\right)^{p-1} x^{r(1-1 / p)} \int_{0}^{x}[f(t)]^{p} t^{p-r-1+r / p} d t \tag{6.50}
\end{align*}
$$

Then by Fubini's theorem,

$$
\begin{aligned}
& \int_{0}^{\infty}\left(\int_{0}^{x} f(t) d t\right)^{p} x^{p-r-1} d x \\
& \quad \leq\left(\frac{p}{r}\right)^{p-1} \int_{0}^{\infty} x^{-1-(r / p)}\left(\int_{0}^{x}[f(t)]^{p} t^{p-r-1+(r / p)} d t\right) d x \\
& \quad=\left(\frac{p}{r}\right)^{p-1} \int_{0}^{\infty}[f(t)]^{p} t^{p-r-1+(r / p)}\left(\int_{t}^{\infty} x^{-1-(r / p)} d x\right) d t \\
& \quad=\left(\frac{p}{r}\right)^{p} \int_{0}^{\infty}[f(t) t]^{p} t^{-r-1} d t
\end{aligned}
$$

The proof of (6.48) proceeds in a similar way.
6.63. LEMMA. If $f \geq 0$ is a non-increasing function on $(0, \infty), 0<p \leq \infty$ and $p_{1} \leq p_{2} \leq \infty$, then

$$
\left(\int_{0}^{\infty}\left[x^{1 / p} f(x)\right]^{p_{2}} \frac{d \lambda(x)}{x}\right)^{1 / p_{2}} \leq C\left(\int_{0}^{\infty}\left[x^{1 / p} f(x)\right]^{p_{1}} \frac{d \lambda(x)}{x}\right)^{1 / p_{1}}
$$

where $C=C\left(p, p_{1}, p_{2}\right)$.

Proof. Since $f$ is non-increasing, we have for any $x>0$

$$
\begin{aligned}
x^{1 / p} f(x) & \leq C\left(\int_{x / 2}^{x}\left[(x / 2)^{1 / p} f(x)\right]^{p_{1}} \frac{d \lambda(y)}{y}\right)^{1 / p_{1}} \\
& \leq C\left(\int_{x / 2}^{x}\left[y^{1 / p} f(x)\right]^{p_{1}} \frac{d \lambda(y)}{y}\right)^{1 / p_{1}} \\
& \leq C\left(\int_{x / 2}^{x}\left[y^{1 / p} f(y)\right]^{p_{1}} \frac{d \lambda(y)}{y}\right)^{1 / p_{1}} \\
& \leq C\left(\int_{0}^{\infty}\left[y^{1 / p} f(y)\right]^{p_{1}} \frac{d \lambda(y)}{y}\right)^{1 / p_{1}}
\end{aligned}
$$

which implies the desired result when $p_{2}=\infty$. The general result follows by writing

$$
\int_{0}^{\infty}\left[x^{1 / p} f(x)\right]^{p_{2}} \frac{d \lambda(x)}{x} \leq \sup _{x>0}\left[x^{1 / p} f(x)\right]^{p_{2}-p_{1}} \int_{0}^{\infty}\left[x^{1 / p} f(x)\right]^{p_{1}} \frac{d \lambda(x)}{x}
$$

6.64. Definition. Let $\mu$ be a nonnegative Radon measure defined on $\mathbb{R}^{n}$ and suppose $f$ is a $\mu$-measurable function defined on $\mathbb{R}^{n}$. Its distribution function, $A_{f}(\cdot)$, is defined by

$$
A_{f}(t):=\mu(\{x:|f(x)|>t\}) .
$$

The non-increasing rearrangement of $f$, denoted by $f^{*}$, is defined as

$$
\begin{equation*}
f^{*}(t)=\inf \left\{\alpha: A_{f}(\alpha) \leq t\right\} \tag{6.51}
\end{equation*}
$$

For example, if $\mu$ is taken as Lebesgue measure, then $f^{*}$ can be identified with that radial function $F$ defined on $\mathbb{R}^{n}$ having the property that, for all $t>0,\{F>t\}$ is a ball centered at the origin whose Lebesgue measure is equal to $\mu(\{x:|f(x)|>$ $t\}$ ). Note that both $f^{*}$ and $A_{f}$ are non-increasing and right continuous. Since $A_{f}$ is right continuous, it follows that the infimum in (6.51) is attained. Therefore if

$$
\begin{equation*}
f^{*}(t)=\alpha, \quad \text { then } A_{f}\left(\alpha^{\prime}\right)>t \text { where } \alpha^{\prime}<\alpha \tag{6.52}
\end{equation*}
$$

Furthermore,

$$
f^{*}(t)>\alpha \quad \text { if and only if } t<A_{f}(\alpha)
$$

Thus, it follows that $\left\{t: f^{*}(t)>\alpha\right\}$ is equal to the interval $\left(0, A_{f}(\alpha)\right)$. Hence, $A_{f}(\alpha)=\lambda\left(\left\{f^{*}>\alpha\right\}\right)$, which implies that $f$ and $f^{*}$ have the same distribution function. Consequently, in view of Theorem 6.60 , observe that for all $1 \leq p \leq \infty$,

$$
\begin{equation*}
\left\|f^{*}\right\|_{p}=\left(\int_{0}^{\infty}\left|f^{*}(t)\right|^{p} d t\right)^{1 / p}=\left(\int_{\mathbb{R}^{n}}|f(x)|^{p} d \mu\right)^{1 / p}=\|f\|_{p ; \mu} \tag{6.53}
\end{equation*}
$$

6.65. Remark. Observe that the $L^{p}$ norm of $f$ relative to the measure $\mu$ is thus expressed as the norm of $f^{*}$ relative to Lebesgue measure.

Notice also that right continuity implies

$$
\begin{equation*}
A_{f}\left(f^{*}(t)\right) \leq t \quad \text { for all } t>0 \tag{6.54}
\end{equation*}
$$

6.66. Lemma. For any $t>0, \sigma>0$, suppose an arbitrary function $f \in L^{p}\left(\mathbb{R}^{n}\right)$ is decomposed as follows: $f=f^{t}+f_{t}$ where

$$
f^{t}(x)= \begin{cases}f(x) & \text { if }|f(x)|>f^{*}\left(t^{\sigma}\right) \\ 0 & i f|f(x)| \leq f^{*}\left(t^{\sigma}\right)\end{cases}
$$

and $f_{t}:=f-f^{t}$. Then

$$
\begin{array}{lr}
\left(f^{t}\right)^{*}(y) \leq f^{*}(y) & \text { if } 0 \leq y \leq t^{\sigma} \\
\left(f^{t}\right)^{*}(y)=0 & \text { if } y>t^{\sigma}  \tag{6.55}\\
\left(f_{t}\right)^{*}(y) \leq f^{*}(y) & \text { if } y>t^{\sigma} \\
\left(f_{t}\right)^{*}(y) \leq f^{*}\left(t^{\sigma}\right) & \text { if } 0 \leq y \leq t^{\sigma}
\end{array}
$$

Proof. We will prove only the first set, since the proof of the other set is similar.

For the first inequality, let $\left[f^{t}\right]^{*}(y)=\alpha$ as in (6.52), and similarly, let $f^{*}(y)=$ $\alpha^{\prime}$. If it were the case that $\alpha^{\prime}<\alpha$, then we would have $A_{f^{t}}\left(\alpha^{\prime}\right)>y$. But, by the definition of $f^{t}$,

$$
\left\{\left|f^{t}\right|>\alpha^{\prime}\right\} \subset\left\{|f|>\alpha^{\prime}\right\}
$$

which would imply

$$
y<A_{f^{t}}\left(\alpha^{\prime}\right) \leq A_{f}\left(\alpha^{\prime}\right) \leq y
$$

a contradiction.
In the second inequality, assume $y>t^{\sigma}$. Now $f^{t}=f \chi_{\left\{|f|>f^{*}\left(t^{\sigma}\right)\right\}}$ and $f^{*}\left(t^{\sigma}\right)=$ $\alpha$ where $A_{f}(\alpha) \leq t^{\sigma}$ as in (6.52). Thus, $A_{f}\left[f^{*}\left(t^{\sigma}\right)\right]=A_{f}(\alpha) \leq t^{\sigma}$ and therefore

$$
\mu\left(\left\{\left|f^{t}\right|>\alpha^{\prime}\right\}\right)=\mu\left(\left\{|f|>f^{*}\left(t^{\sigma}\right)\right\}\right) \leq t^{\sigma}<y
$$

for all $\alpha^{\prime}>0$. This implies $\left(f^{t}\right)^{*}(y)=0$.
6.67. Definition. Suppose $(X \mathcal{M}, \mu)$ is a measure space and let $(p, q)$ be a pair of numbers such that $1 \leq p, q<\infty$. Also, let $\mu$ be a Radon measure defined on $X$ and suppose $T$ is an sub-additive operator defined on $L^{p}(X)$ whose values are $\mu$-measurable functions. Thus, $T(f)$ is a $\mu$-measurable function on $X$ and we will write $T f:=T(f)$. The operator $T$ is said to be of weak-type $(p, q)$ if there is a constant $C$ such that for any $f \in L^{p}(X, \mu)$ and $\alpha>0$,

$$
\mu(\{x:|(T f)(x)|>\alpha\}) \leq\left(\alpha^{-1} C\|f\|_{p ; \mu}\right)^{q}
$$

$T$ is said to be of strong type $(p, q)$ if there is a constant $C$ such that $\|T f\|_{q ; \mu} \leq$ $C\|f\|_{p ; \mu}$ for all $f \in L^{p}(X, \mu)$.
6.68. Theorem (Marcinkiewicz Interpolation Theorem). Let ( $p_{0}, q_{0}$ ) and ( $p_{1}, q_{1}$ ) be pairs of numbers such that $1 \leq p_{i} \leq q_{i}<\infty, i=0,1$, and $q_{0} \neq q_{1}$. Let $\mu$ be a Radon measure defined on $\mathbb{R}^{n}$ and suppose $T$ is an sub-additive operator defined on $L^{p_{0}}\left(\mathbb{R}^{n}\right)+L^{p_{1}}\left(\mathbb{R}^{n}\right)$ whose values are $\mu$-measurable functions. Suppose $T$ is simultaneously of weak-types $\left(p_{0}, q_{0}\right)$ and $\left(p_{1}, q_{1}\right)$. If $0<\theta<1$, and

$$
\begin{align*}
& 1 / p=\frac{1-\theta}{p_{0}}+\frac{\theta}{p_{1}} \\
& 1 / q=\frac{1-\theta}{q_{0}}+\frac{\theta}{q_{1}} \tag{6.56}
\end{align*}
$$

then $T$ is of strong type $(p, q)$; that is,

$$
\|T f\|_{q ; \mu} \leq C\|f\|_{p ; \mu}, \quad f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

where $C=C\left(p_{0}, q_{0}, p_{1}, q_{1}, \theta\right)$.

Proof. The easiest case arises when $p_{0}=p_{1}$ and is left as an exercise.
Henceforth, assume $p_{0}<p_{1}$. Let $(T f)^{*}(t)=\alpha$ as in (6.52). Then for $\alpha^{\prime}<$ $\alpha, A_{T f}\left(\alpha^{\prime}\right)>t$. The weak-type $\left(p_{0}, q_{0}\right)$ assumption on $T$ implies

$$
\begin{aligned}
\alpha^{\prime} & \leq C_{0}\left(A_{T f}\left(\alpha^{\prime}\right)\right)^{-1 / q_{0}}\|f\|_{p_{0} ; \mu} \\
& <C_{0} t^{-1 / q_{0}}\|f\|_{p_{0} ; \mu}
\end{aligned}
$$

whenever $f \in L^{p_{0}}\left(\mathbb{R}^{n}\right)$. Since $\alpha^{\prime}<C_{0} t^{-1 / q_{0}}\|f\|_{p_{0} ; \mu}$ for all $\alpha^{\prime}<\alpha=(T f)^{*}(t)$, it follows that

$$
\begin{equation*}
(T f)^{*}(t) \leq C_{0} t^{-1 / q_{0}}\|f\|_{p_{0} ; \mu} \tag{6.57}
\end{equation*}
$$

A similarly argument shows that if $f \in L^{p_{1}}\left(\mathbb{R}^{n}\right)$, then

$$
\begin{equation*}
(T f)^{*}(t) \leq C_{1} t^{-1 / q_{1}}\|f\|_{p_{1} ; \mu} \tag{6.58}
\end{equation*}
$$

We now appeal to Lemma 6.66 where $\sigma$ is taken as

$$
\begin{equation*}
\sigma:=\frac{1 / q_{0}-1 / q}{1 / p_{0}-1 / p}=\frac{1 / q-1 / q_{1}}{1 / p-1 / p_{1}} \tag{6.59}
\end{equation*}
$$

Recall the decomposition $f=f^{t}+f_{t}$; since $p_{0}<p<p_{1}$, observe from (6.55) that $f^{t} \in L^{p_{0}}\left(\mathbb{R}^{n}\right)$ and $f_{t} \in L^{p_{1}}\left(\mathbb{R}^{n}\right)$. Also, we leave the following as an exercise 6.3:

$$
\begin{equation*}
(T f)^{*}(t) \leq\left(T f^{t}\right)^{*}(t / 2)+\left(T f_{t}\right)^{*}(t / 2) \tag{6.60}
\end{equation*}
$$

Since $p_{i} \leq q_{i}, i=0,1$, by (6.56) we have $p \leq q$. Thus, we obtain

$$
\begin{align*}
\left\|(T f)^{*}\right\|_{q}= & \left(\int_{0}^{\infty}\left(t^{1 / q}(T f)^{*}(t)\right)^{q} \frac{d t}{t}\right)^{1 / q} \\
\leq & C\left(\int_{0}^{\infty}\left(t^{1 / q}(T f)^{*}(t)\right)^{p} \frac{d t}{t}\right)^{1 / p} \quad \text { by Lemma } 6.63 \\
\leq & C\left(\int_{0}^{\infty}\left(t^{1 / q}\left(T f^{t}\right)^{*}(t)\right)^{p} \frac{d t}{t}\right)^{1 / p} \\
& +C\left(\int_{0}^{\infty}\left(t^{1 / q}\left(T f_{t}\right)^{*}(t)\right)^{p} \frac{d t}{t}\right)^{1 / p} \quad \text { by }(6.60)  \tag{6.61}\\
\leq & C\left(\int_{0}^{1}\left(t^{1 / q-1 / q_{0}}\left\|f^{t}\right\|_{p_{0}}\right)^{p} \frac{d t}{t}\right)^{1 / p} \quad \text { by }(6.57)  \tag{6.62}\\
& +C\left(\int_{0}^{1}\left(t^{1 / q-1 / q_{1}}\left\|f_{t}\right\|_{p_{1}}\right)^{p} \frac{d t}{t}\right)^{1 / p} \quad \text { by }(6.58) \tag{6.63}
\end{align*}
$$

With $\sigma$ defined by (6.59) we estimate the last two integrals with an appeal to Lemma 6.63 and write

$$
\begin{aligned}
\left\|f^{t}\right\|_{p_{0}} & =\left(\int_{0}^{\infty}\left(y^{1 / p_{0}}\left(f^{t}\right)^{*}(y)\right)^{p_{0}} \frac{d \lambda(y)}{y}\right)^{1 / p_{0}} \\
& \leq C \int_{0}^{t^{\sigma}} y^{1 / p_{0}}\left(f^{t}\right)^{*}(y) \frac{d \lambda(y)}{y} \\
& \leq C \int_{0}^{t^{\sigma}} y^{1 / p_{0}} f^{*}(y) \frac{d \lambda(y)}{y}
\end{aligned}
$$

Inserting this estimate for $\left\|f^{t}\right\|_{p}$ into (6.62) we obtain

$$
\begin{aligned}
& \left(\int_{0}^{1}\left(t^{1 / q-1 / q_{0}}\left\|f^{t}\right\|_{p_{0}}\right)^{p} \frac{d t}{t}\right)^{1 / p} \\
& \quad \leq C\left(\int_{0}^{1}\left(\left(t^{1 / q-1 / q_{0}} \int_{0}^{t^{\sigma}} y^{1 / p_{0}} f^{*}(y) \frac{d \lambda(y)}{y}\right)^{p} \frac{d t}{t}\right)^{1 / p}\right.
\end{aligned}
$$

Thus, to estimate (6.62) we analyze

$$
\int_{0}^{\infty}\left(t^{1 / q-1 / q_{0}} \int_{0}^{t^{\sigma}} y^{1 / p_{0}} f^{*}(y) \frac{d \lambda(y)}{y}\right)^{p} \frac{d t}{t}
$$

which, under the change of variables $t^{\sigma} \mapsto s$, becomes

$$
\frac{1}{\sigma} \int_{0}^{\infty}\left(s^{1 / p-1 / p_{0}} \int_{0}^{s} y^{1 / p_{0}-1} f^{*}(y) d \lambda(y)\right)^{p} \frac{d s}{s}
$$

which is equal to

$$
\frac{1}{\sigma} \int_{0}^{\infty}\left(s^{1 / p-1 / p_{0}-1 / p} \int_{0}^{s} y^{1 / p_{0}-1} f^{*}(y) d \lambda(y)\right)^{p} d s
$$

Now apply Hardy's inequality (6.47) with $-r-1=-p / p_{0}$ and $f(y)=y^{1 / p_{0}-1} f^{*}(y)$ to obtain

$$
\begin{aligned}
\left(\int_{0}^{1}\left(t^{1 / q-1 / q_{0}}\left\|f^{t}\right\|_{p_{0} ; \mu}\right)^{p} \frac{d t}{t}\right)^{1 / p} & \leq C(p, r)\left(\int_{0}^{\infty} f(y)^{p} y^{p-r-1}\right)^{1 / p} \\
& =C(p, r)\left\|f^{*}\right\|_{p}
\end{aligned}
$$

The estimate of

$$
\int_{0}^{\infty}\left(t^{1 / q-1 / q_{1}} \int_{0}^{t^{\sigma}} y^{1 / p_{1}} f^{*}(y) \frac{d y}{y}\right)^{p} \frac{d t}{t}
$$

proceeds in a similar way, and thus our result is established.

## Exercises for Section 6.13

1. Let $f$ be a measurable function on a measure space $(X, \mathcal{M})$. Define $A_{f}(s)=$ $\mu(\{|f|>s\})$. The nonincreasing rearrangement of $f$ on $(0, \infty)$ is defined as

$$
f^{*}(t):=\inf \left\{s: A_{f}(s) \leq t\right\}
$$

Prove the following:
(i) $f^{*}$ is continuous from the right.
(ii) $A_{f^{*}}(s)=A_{f}(s)$ for all $s$ in case $\mu$ is Lebesgue measure on $\mathbb{R}^{n}$.
2. Prove the Marcinkiewicz Interpolation Theorem in the case when $p_{0}=p_{1}$.
3. If $f=f_{1}+f_{2}$ prove that

$$
\begin{equation*}
(T f)^{*}(t) \leq\left(T f_{1}\right)^{*}(t / 2)+\left(T f_{2}\right)^{*}(t / 2) \tag{6.64}
\end{equation*}
$$

## CHAPTER 7

## Differentiation

### 7.1. Covering Theorems

> Certain covering theorems, such as the Vitali Covering Theorem, will be developed in this section. These covering theorems are of essential importance in the theory of differentiation of measures.

We depart from the theory of abstract measure spaces encountered in previous chapters and focus on certain aspects of functions defined in $\mathbb{R}$. A major result in elementary analysis is the fundamental theorem of calculus, which states that a $C^{1}$ function can be expressed as the integral of its derivative. One of the main objectives of this chapter is to show that this result still holds for a more general class of functions. In fact, we will determine precisely those functions for which the fundamental theorem holds. We will take a broader view of differentiation by developing a framework for differentiation of measures. This will include the usual notion of differentiability of a function. The following result, whose proof is left as an exercise, will serve to motivate our point of view.
7.1. Remark. Suppose $\mu$ is a Borel measure on $\mathbb{R}$ and let

$$
F(x)=\mu((-\infty, x]) \quad \text { for } \quad x \in \mathbb{R}
$$

Then the following two statements are equivalent:
(i) $F$ is differentiable at $x_{0}$ and $F^{\prime}\left(x_{0}\right)=c$
(ii) For every $\varepsilon>0$ there exists $\delta>0$ such that

$$
\left|\frac{\mu(I)}{\lambda(I)}-c\right|<\varepsilon
$$

whenever $I$ is a half-open interval whose left or right endpoint is $x_{0}$ and $\lambda(I)<\delta$.
Condition (ii) may be interpreted as the derivative of $\mu$ with respect to Lebesgue measure, $\lambda$. This concept will be developed more fully throughout this chapter. As an example in this framework, let $f \in L^{1}(\mathbb{R})$ be nonnegative and define a measure $\mu$ by

$$
\begin{equation*}
\mu(E)=\int_{E} f(y) d \lambda(y) \tag{7.1}
\end{equation*}
$$

for every Lebesgue measurable set $E$. Then the function $F$ introduced above can be expressed as

$$
F(x)=\int_{-\infty}^{x} f(y) d \lambda(y)
$$

Of course, the derivative of $F$ at $x_{0}$ is the limit

$$
\begin{equation*}
\lim _{h \rightarrow 0} \frac{1}{h} \int_{x_{0}}^{x_{0}+h} f(y) d \lambda(y) \tag{7.2}
\end{equation*}
$$

This, in turn, is equivalent to statement (ii) above. Given $f \in L^{1}\left(\mathbb{R}^{n}\right)$, we define $\mu$ as in (7.1) and we consider

$$
\begin{equation*}
\lim _{r \rightarrow 0} \frac{\mu\left[B\left(x_{0}, r\right)\right]}{\lambda\left[B\left(x_{0}, r\right)\right]}=\lim _{r \rightarrow 0} \frac{1}{\lambda\left[B\left(x_{0}, r\right)\right]} \int_{B\left(x_{0}, r\right)} f(y) d \lambda(y) \tag{7.3}
\end{equation*}
$$

At this stage we know nothing about the existence of the limit.
7.2. Remark. Strictly speaking, (7.3) is not the precise analog of (7.2) since we have considered only balls $B\left(x_{0}, r\right)$ centered at $x_{0}$. In $\mathbb{R}$ this would exclude the use of intervals whose left endpoint is $x_{0}$ as required by (7.2). Nevertheless, in our development we choose to use the family of open concentric balls for several reasons. First, they are slightly easier to employ than nonconcentric balls; second, we will see that it is immaterial to the main results of the theory whether or not concentric balls are used (see Theorem 7.17). Finally, in the development of the derivative of $\mu$ relative to an arbitrary measure $\nu$, it is important that concentric balls are used. Thus, we formally introduce the notation

$$
\begin{equation*}
D_{\lambda} \mu\left(x_{0}\right)=\lim _{r \rightarrow 0} \frac{\mu\left[B\left(x_{0}, r\right)\right]}{\lambda\left[B\left(x_{0}, r\right)\right]} \tag{7.4}
\end{equation*}
$$

which is the derivative of $\mu$ with respect to $\lambda$.
One of the major objectives of the next section is to prove that the limit in (7.4) exists $\lambda$ almost everywhere and to see how it relates to the results surrounding the Radon-Nikodym Theorem. In particular, in view of Theorems 4.32 and (7.1), it will follow from our development that a nondecreasing function is differentiable almost everywhere. In the following we will use the following notation: Given $B=B(x, r)$ we will call $B(x, 5 r)$ the enlargement of $B$ and denote it by $\widehat{B}$. The next lemma states that any collection of balls (that may be either open or closed) whose radii are bounded has a countable disjoint subcollection with the property that the union of their enlargements contains the union of the original collection. We emphasize here that the point of the lemma is that the subcollection consists of disjoint elements, a very important consideration since countable additivity plays a central role in measure theory.
7.3. Theorem. Let $\mathcal{G}$ be a family of closed or open balls in $\mathbb{R}^{n}$ with

$$
R:=\sup \{\operatorname{diam} B: B \in \mathcal{G}\}<\infty
$$

Then there is a countable subfamily $\mathcal{F} \subset \mathcal{G}$ of pairwise disjoint elements such that

$$
\bigcup\{B: B \in \mathcal{G}\} \subset \bigcup\{\widehat{B}: B \in \mathcal{F}\}
$$

In fact, for each $B \in \mathcal{G}$ there exists $B^{\prime} \in \mathcal{F}$ such that $B \cap B^{\prime} \neq 0$ and $B \subset \widehat{B}^{\prime}$.
Proof. Throughout this proof, we will adopt the following notation: If $A$ is set and $\mathcal{F}$ a family of sets, then we use the notation $A \cap \mathcal{F} \neq \emptyset$ to mean that $A \cap B \neq \emptyset$ for some $B \in \mathcal{F}$.

Let $a$ be a number such that $1<a<a^{1+2 R}<2$. For $j=1,2, \ldots$ let

$$
\mathcal{G}_{j}=\left\{B:=B(x, r) \in \mathcal{G}: a^{|x|-j}<\frac{r}{R} \leq a^{|x|-j+1}\right\}
$$

and observe that $\mathcal{G}=\bigcup_{j=1}^{\infty} \mathcal{G}_{j}$. Since $r / R<1$ and $a>1$, observe that for any $B=B(x, r) \in \mathcal{G}_{j}$ we have

$$
|x|<j \text { and } r>a^{-j} R .
$$

Hence the elements of $\mathcal{G}_{j}$ are centered at points $x \in B(0, j)$ and their radii are bounded away from zero; this implies that there is a number $M_{j}>0$ depending only on $a, j$ and $R$ such that any disjoint subfamily of $\mathcal{G}_{j}$ has at most $M_{j}$ elements.

The family $\mathcal{F}$ will be of the form

$$
\mathcal{F}=\bigcup_{j=0}^{\infty} \mathcal{F}_{j}
$$

where the $\mathcal{F}_{j}$ are finite, disjoint families defined inductively as follows. We set $\mathcal{F}_{0}=\emptyset$. Let $\mathcal{F}_{1}$ be the largest (in the sense of inclusion) disjoint subfamily of $\mathcal{G}_{1}$. Note that $\mathcal{F}_{1}$ can have no more than $M_{1}$ elements. Proceeding by induction, we assume that $\mathcal{F}_{j-1}$ has been determined, and then define $\mathcal{H}_{j}$ as the largest disjoint subfamily of $\mathcal{G}_{j}$ with the property that $B \cap \mathcal{F}_{j-1}=\emptyset$ for each $B \in \mathcal{H}_{j}$. Note that the number of elements in $\mathcal{H}_{j}$ could be 0 but no more than $M_{j}$. Define

$$
\mathcal{F}_{j}:=\mathcal{F}_{j-1} \cup \mathcal{H}_{j}
$$

We claim that the family $\mathcal{F}:=\bigcup_{j=1}^{\infty} \mathcal{F}_{j}$ has the required properties: that is, we will show that

$$
\begin{equation*}
B \subset \bigcup\{\widehat{B}: B \in \mathcal{F}\} \quad \text { for each } B \in \mathcal{G} \tag{7.5}
\end{equation*}
$$

To verify this, first note that $\mathcal{F}$ is a disjoint family. Next, select $B:=B(x, r) \in$ $\mathcal{G}$ which implies $B \in \mathcal{G}_{j}$ for some $j$. If $B \cap \mathcal{H}_{j} \neq \emptyset$, then there exists $B^{\prime}:=B\left(x^{\prime}, r^{\prime}\right) \in$ $\mathcal{H}_{j}$ such that $B \cap B^{\prime} \neq \emptyset$, in which case

$$
\begin{equation*}
\frac{r^{\prime}}{R} \geq a^{\left|x^{\prime}\right|-j} \tag{7.6}
\end{equation*}
$$

On the other hand, if $B \cap \mathcal{H}_{j}=\emptyset$, then $B \cap \mathcal{F}_{j-1} \neq \emptyset$, for otherwise the maximality of $\mathcal{H}_{j}$ would be violated. Thus, there exists $B^{\prime}=B\left(x^{\prime}, r^{\prime}\right) \in \mathcal{F}_{j-1}$ such that $B \cap B^{\prime} \neq \emptyset$ and in this case

$$
\begin{equation*}
\frac{r^{\prime}}{R}>a^{\left|x^{\prime}\right|-j+1}>a^{\left|x^{\prime}\right|-j} \tag{7.7}
\end{equation*}
$$

Since

$$
\frac{r}{R} \leq a^{|x|-j+1}
$$

it follows from (7.6) and (7.7) that

$$
r \leq a^{|x|-j+1} R \leq a^{\left(|x|-\left|x^{\prime}\right|+1\right)} r^{\prime}
$$

Since $a$ was chosen so that $1<a<a^{1+2 R}<2$, we have

$$
r \leq a^{1+|x|-\left|x^{\prime}\right|} r^{\prime} \leq a^{1+\left|x-x^{\prime}\right|} r^{\prime} \leq a^{1+2 R} r^{\prime} \leq 2 r^{\prime}
$$

This implies that $B \subset \widehat{B}^{\prime}$, because if $z \in B(x, r)$ and $y \in B \cap B^{\prime}$, then

$$
\begin{aligned}
\left|z-x^{\prime}\right| & \leq|z-x|+|x-y|+\left|y-x^{\prime}\right| \\
& \leq r+r+r^{\prime} \\
& \leq 5 r^{\prime}
\end{aligned}
$$

If we assume a bit more about $\mathcal{G}$ we can show that the union of elements in $\mathcal{F}$ contains almost all of the union $\bigcup\{B: B \in \mathcal{G}\}$. This requires the following definition.
7.4. Definition. A collection $\mathcal{G}$ of balls is said to cover a set $E \subset \mathbb{R}^{n}$ in the sense of Vitali if for each $x \in E$ and each $\varepsilon>0$, there exists $B \in \mathcal{G}$ containing $x$ whose radius is positive and less than $\varepsilon$. We also say that $\mathcal{G}$ is a Vitali covering of $E$. Note that if $\mathcal{G}$ is a Vitali covering of a set $E \subset \mathbb{R}^{n}$ and $R>0$ is arbitrary, then $\mathcal{G} \bigcap\{B: \operatorname{diam} B<R\}$ is also a Vitali covering of $E$.
7.5. Theorem. Let $\mathcal{G}$ be a family of closed balls that covers a set $E \subset \mathbb{R}^{n}$ in the sense of Vitali. Then with $\mathcal{F}$ as in Theorem 7.3, we have

$$
E \backslash \bigcup\left\{B: B \in \mathcal{F}^{*}\right\} \subset \bigcup\left\{\widehat{B}: B \in \mathcal{F} \backslash \mathcal{F}^{*}\right\}
$$

for each finite collection $\mathcal{F}^{*} \subset \mathcal{F}$.

Proof. Since $\mathcal{G}$ is a Vitali covering of $E$, there is no loss of generality if we assume that the radius of each ball in $\mathcal{G}$ is less than some fixed number $R$. Let $\mathcal{F}$ be as in Theorem 7.3 and let $\mathcal{F}^{*}$ be any finite subfamily of $\mathcal{F}$. Since $R^{n} \backslash \cup\left\{B: B \in \mathcal{F}^{*}\right\}$ is open, for each $x \in E \backslash \cup\left\{B: B \in \mathcal{F}^{*}\right\}$ there exists $B \in \mathcal{G}$ such that $x \in B$ and $B \cap\left[\cup\left\{B: B \in \mathcal{F}^{*}\right\}\right]=\emptyset$. From Theorem 7.3, there is $B_{1} \in \mathcal{F}$ such that $B \cap B_{1} \neq \emptyset$ and $\widehat{B}_{1} \supset B$. Since $\mathcal{F}^{*}$ is disjoint, it follows that $B_{1} \notin \mathcal{F}^{*}$ since $B \cap B_{1} \neq \emptyset$. Therefore

$$
x \in \widehat{B}_{1} \subset \bigcup\left\{\widehat{B}: B \in \mathcal{F} \backslash \mathcal{F}^{*}\right\}
$$

7.6. Remark. The preceding result and the next one are not needed in the sequel, although they are needed in some of the Exercises, such as Exercise 3, Section 7.9. We include them because they are frequently used in the analysis literature and because they follow so easily from the main result, Theorem 7.3.

Theorem 7.5 states that any finite family $\mathcal{F}^{*} \subset \mathcal{F}$ along with the enlargements of $\mathcal{F} \backslash \mathcal{F}^{*}$ provide a covering of $E$. But what covering properties does $\mathcal{F}$ itself have? The next result shows that $\mathcal{F}$ covers almost all of $E$.
7.7. Theorem. Let $\mathcal{G}$ be a family of closed balls that covers a (possibly nonmeasurable) set $E \subset \mathbb{R}^{n}$ in the sense of Vitali. Then there exists a countable disjoint subfamily $\mathcal{F} \subset \mathcal{G}$ such that

$$
\lambda(E \backslash \bigcup\{B: B \in \mathcal{F}\})=0
$$

Proof. First, assume that $E$ is a bounded set. Then we may as well assume that each ball in $\mathcal{G}$ is contained in some bounded open set $H \supset E$. Let $\mathcal{F}$ be the subfamily of disjoint balls provided by Theorem 7.3 and Corollary 7.5. Since all elements of $\mathcal{F}$ are disjoint and contained in the bounded set $H$, we have

$$
\begin{equation*}
\sum_{B \in \mathcal{F}} \lambda(B) \leq \lambda(H)<\infty \tag{7.8}
\end{equation*}
$$

Now, by Corollary 7.5 , for any finite subfamily $\mathcal{F}^{*} \subset \mathcal{F}$, we obtain

$$
\begin{aligned}
\lambda^{*}(E \backslash \cup\{B: B \in \mathcal{F}\}) & \leq \lambda^{*}\left(E \backslash \cup\left\{B: B \in \mathcal{F}^{*}\right\}\right) \\
& \leq \lambda^{*}\left(\cup\left\{\widehat{B}: B \in \mathcal{F} \backslash \mathcal{F}^{*}\right)\right. \\
& \leq \sum_{B \in \mathcal{F} \backslash \mathcal{F}^{*}} \lambda(\widehat{B}) \\
& \leq 5^{n} \sum_{B \in \mathcal{F} \backslash \mathcal{F}^{*}} \lambda(B) .
\end{aligned}
$$

Referring to (7.8), we see that the last term can be made arbitrarily small by an appropriate choice of $\mathcal{F}^{*}$. This establishes our result in case $E$ is bounded.

The general case can be handled by observing that there is a countable family $\left\{C_{k}\right\}_{k=1}^{\infty}$ of disjoint open cubes $C_{k}$ such that

$$
\lambda\left(\mathbb{R}^{n} \backslash \bigcup_{k=1}^{\infty} C_{k}\right)=0
$$

The details are left to the reader.

## Exercises for Section 7.1

1. (a) Prove for each open set $U \subset \mathbb{R}^{n}$ there exists a collection $\mathcal{F}$ of disjoint closed balls contained in $U$ such that

$$
\lambda(U \backslash \bigcup\{B: B \in \mathcal{F}\})=0
$$

(b) Thus, $U=\bigcup_{B \in \mathcal{F}} B \cup N$ where $\lambda(N)=0$. Prove that $N \neq \emptyset$. (Hint: to show that $N \neq \emptyset$, consider the proof in $\mathbb{R}^{2}$. Then consider the intersection of $U$ with a line. This intersection is an open subset of the line and note that the closed balls become closed intervals.)
2. Let $\mu$ be a finite Borel measure on $\mathbb{R}^{n}$ with the property that $\mu[B(x, 2 r)] \leq$ $c \mu[B(x, r)]$ for all $x \in \mathbb{R}^{n}$ and all $0<r<\infty$, where $c$ is a constant independent of $x$ and $r$. Prove that the Vitali covering theorem, Theorem 7.7, is valid with $\lambda$ replaced by $\mu$.
3. Supply the proof of Theorem 7.1.
4. Prove the following alternate version of Theorem 7.7. Let $E \subset \mathbb{R}^{n}$ be an arbitrary set, possibly nonmeasurable. Suppose $\mathcal{G}$ is a family of closed cubes with the property that for each $x \in E$ and each $\varepsilon>0$ there exists a cube $C \in \mathcal{G}$ containing $x$ whose diameter is less than $\varepsilon$. Prove that there exists a countable disjoint subfamily $\mathcal{F} \subset \mathcal{G}$ such that

$$
\lambda(E-\bigcup\{C: C \in \mathcal{F}\})=0
$$

5. With the help of the preceding exercise, prove that for each open set $U \subset \mathbb{R}^{n}$, there exists a countable family $\mathcal{F}$ of closed, disjoint cubes, each contained in $U$ such that

$$
\lambda(U \backslash \bigcup\{C: C \in \mathcal{F}\})=0
$$

### 7.2. Lebesgue Points

[^1]7.8. Definition. With each $f \in L^{1}\left(\mathbb{R}^{n}\right)$, we associate its maximal function, $\boldsymbol{M} f$, which is defined as
$$
\boldsymbol{M} f(x):=\sup _{r>0} f_{B(x, r)}|f| d \lambda
$$
where
$$
f_{E}|f| d \lambda:=\frac{1}{\lambda[E]} \int_{E}|f| d \lambda
$$
denotes the integral average of $|f|$ over an arbitrary measurable set $E$. In other words, $\boldsymbol{M} f(x)$ is the upper envelope of integral averages of $|f|$ over balls centered at $x$.

Clearly, $\boldsymbol{M} f: \mathbb{R}^{n} \rightarrow \overline{\mathbb{R}}$ is a nonnegative function. Furthermore, it is Lebesgue measurable. To see this, note that for each fixed $r>0$,

$$
x \mapsto f_{B(x, r)}|f| d \lambda
$$

is a continuous function of $x$ (see Exercise 4, Section 7.2). Therefore, we see that $\{\boldsymbol{M} f>t\}$ is an open set for each real number $t$, thus showing that $\boldsymbol{M} f$ is lower semicontinuous and therefore measurable.

The next question is whether $M f$ is integrable over $\mathbb{R}^{n}$. In order for this to be true, it follows from Theorem 6.59 that it would be necessary that

$$
\begin{equation*}
\int_{0}^{\infty} \lambda(\{\boldsymbol{M} f>t\}) d \lambda(t)<\infty \tag{7.9}
\end{equation*}
$$

It turns out that $\boldsymbol{M} f$ is never integrable unless $f$ is identically zero (see Exercise 3, Section 7.2). However, the next result provides an estimate of how the measure of the set $\{\boldsymbol{M} f>t\}$ becomes small as $t$ increases. It also shows that inequality (7.9) fails to be true by only a small margin.
7.9. Theorem (Hardy-Littlewood). If $f \in L^{1}\left(\mathbb{R}^{n}\right)$, then

$$
\lambda[\{\boldsymbol{M} f>t\}] \leq \frac{5^{n}}{t} \int_{\mathbb{R}^{n}}|f| d \lambda
$$

for every $t>0$.
Proof. For fixed $t>0$, the definition implies that for each $x \in\{M f>t\}$ there exists a ball $B_{x}$ centered at $x$ such that

$$
f_{B_{x}}|f| d \lambda>t
$$

or what is the same

$$
\begin{equation*}
\frac{1}{t} \int_{B_{x}}|f| d \lambda>\lambda\left(B_{x}\right) \tag{7.10}
\end{equation*}
$$

Since $f$ is integrable and $t$ is fixed, the radii of all balls satisfying (7.10) is bounded. Thus, with $\mathcal{G}$ denoting the family of these balls, we may appeal to Lemma 7.3 to obtain a countable subfamily $\mathcal{F} \subset \mathcal{G}$ of disjoint balls such that

$$
\{\boldsymbol{M} f>t\} \subset \bigcup\{\widehat{B}: B \in \mathcal{F}\}
$$

Therefore,

$$
\begin{aligned}
\lambda(\{\boldsymbol{M} f>t\}) & \leq \lambda\left(\bigcup_{B \in \mathcal{F}} \widehat{B}\right) \\
& \leq \sum_{B \in \mathcal{F}} \lambda(\widehat{B}) \\
& =5^{n} \sum_{B \in \mathcal{F}} \lambda(B) \\
& <\frac{5^{n}}{t} \sum_{B \in \mathcal{F}} \int_{B}|f| d \lambda \\
& \leq \frac{5^{n}}{t} \int_{\mathbb{R}^{n}}|f| d \lambda
\end{aligned}
$$

which establishes the desired result.

We now appeal to the results of Section 6.13 concerning the Marcinkiewicz Interpolation Theorem. Clearly, the operator $\boldsymbol{M}$ is sub-additive and our previous result shows that it is of weak type $(1,1)$, (see Definition 6.67). Also, it is clear that

$$
\|\boldsymbol{M} f\|_{\infty} \leq\|f\|_{\infty}
$$

for all $f \in L^{\infty}$. Therefore, we appeal to the Marcinkiewicz Interpolation Theorem to conclude that $M$ is of strong type $(p, p)$. That is, we have
7.10. Corollary. There exists a constant $C>0$ such that

$$
\|\boldsymbol{M} f\|_{p} \leq C p(p-1)^{-1}\|f\|_{p}
$$

whenever $1<p<\infty$ and $f \in L^{p}\left(\mathbb{R}^{n}\right)$.
If $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$ is continuous, it follows from elementary considerations that

$$
\begin{equation*}
\lim _{r \rightarrow 0} f_{B(x, r)} f(y) d \lambda(y)=f(x) \quad \text { for } \quad x \in \mathbb{R}^{n} \tag{7.11}
\end{equation*}
$$

Since Lusin's Theorem tells us that a measurable function is almost continuous, one might suspect that (7.11) is true in some sense for an integrable function. Indeed, we have the following.
7.11. Theorem. If $f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$, then

$$
\begin{equation*}
\lim _{r \rightarrow 0} f_{B(x, r)} f(y) d \lambda(y)=f(x) \tag{7.12}
\end{equation*}
$$

for a.e. $x \in \mathbb{R}^{n}$.
Proof. Since the limit in (7.12) depends only on the values of $f$ in an arbitrarily small neighborhood of $x$, and since $\mathbb{R}^{n}$ is a countable union of bounded measurable sets, we may assume without loss of generality that $f$ vanishes on the complement of a bounded set. Choose $\varepsilon>0$. From Exercise 12, Section 6.5, we can find a continuous function $g \in L^{1}\left(\mathbb{R}^{n}\right)$ such that

$$
\int_{\mathbb{R}^{n}}|f(y)-g(y)| d \lambda(y)<\varepsilon
$$

For each such $g$ we have

$$
\lim _{r \rightarrow 0} f_{B(x, r)} g(y) d \lambda(y)=g(x)
$$

for every $x \in \mathbb{R}^{n}$. This implies

$$
\begin{align*}
& \limsup _{r \rightarrow 0}\left|f_{B(x, r)} f(y) d \lambda(y)-f(x)\right| \\
& =\limsup _{r \rightarrow 0} \mid f_{B(x, r)}[f(y)-g(y)] d \lambda(y)  \tag{7.13}\\
& \quad+\left(f_{B(x, r)} g(y) d \lambda(y)-g(x)\right)+[g(x)-f(x)] \mid \\
& \leq M(f-g)(x)+0+|f(x)-g(x)|
\end{align*}
$$

For each positive number $t$ let

$$
\begin{gathered}
E_{t}=\left\{x: \limsup _{r \rightarrow 0}\left|f_{B(x, r)} f(y) d \lambda(y)-f(x)\right|>t\right\} \\
F_{t}=\{x:|f(x)-g(x)|>t\}
\end{gathered}
$$

and

$$
H_{t}=\{x: \boldsymbol{M}(f-g)(x)>t\} .
$$

Then, by (7.13), $E_{t} \subset F_{t / 2} \cup H_{t / 2}$. Furthermore,

$$
t \lambda\left(F_{t}\right) \leq \int_{F_{t}}|f(y)-g(y)| d \lambda(y)<\varepsilon
$$

and Theorem 7.9 implies

$$
\lambda\left(H_{t}\right) \leq \frac{5^{n} \varepsilon}{t}
$$

Hence

$$
\lambda\left(E_{t}\right) \leq 2 \frac{\varepsilon}{t}+2 \frac{5^{n} \varepsilon}{t}
$$

Since $\varepsilon$ is arbitrary, we conclude that $\lambda\left(E_{t}\right)=0$ for all $t>0$, thus establishing the conclusion.

The theorem states that

$$
\begin{equation*}
\lim _{r \rightarrow 0} f_{B(x, r)} f(y) d \lambda(y) \tag{7.14}
\end{equation*}
$$

exists for a.e. $x$ and that the limit defines a function that is equal to $f$ almost everywhere. The limit in (7.14) provides a way to define the value of $f$ at $x$ that is independent of the choice of representative in the equivalence class of $f$. Observe that (7.12) can be written as

$$
\lim _{r \rightarrow 0} f_{B(x, r)}[f(y)-f(x)] d \lambda(y)=0
$$

It is rather surprising that Theorem 7.11 implies the following apparently stronger result.
7.12. Theorem. If $f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$, then

$$
\begin{equation*}
\lim _{r \rightarrow 0} f_{B(x, r)}|f(y)-f(x)| d \lambda(y)=0 \tag{7.15}
\end{equation*}
$$

for a.e. $x \in \mathbb{R}^{n}$.

Proof. For each rational number $\rho$ apply Theorem 7.11 to conclude that there is a set $E_{\rho}$ of measure zero such that

$$
\begin{equation*}
\lim _{r \rightarrow 0} f_{B(x, r)}|f(y)-\rho| d \lambda(y)=|f(x)-\rho| \tag{7.16}
\end{equation*}
$$

for all $x \notin E_{\rho}$. Thus, with

$$
E:=\bigcup_{\rho \in \mathbb{Q}} E_{\rho}
$$

we have $\lambda(E)=0$. Moreover, for $x \notin E$ and $\rho \in \mathbb{Q}$, then, since $|f(y)-f(x)|<$ $|f(y)-\rho|+|f(x)-\rho|,(7.16)$ implies

$$
\limsup _{r \rightarrow 0} f_{B(x, r)}|f(y)-f(x)| d \lambda(y) \leq 2|f(x)-\rho|
$$

Since

$$
\inf \{|f(x)-\rho|: \rho \in \mathbb{Q}\}=0
$$

the proof is complete.

A point $x$ for which (7.15) holds is called a Lebesgue point of $f$. Thus, almost all points are Lebesgue points for any $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$.

An important special case of Theorem 7.11 is when $f$ is taken as the characteristic function of a set. For $E \subset \mathbb{R}^{n}$ a Lebesgue measurable set, let

$$
\begin{align*}
& \bar{D}(E, x)=\limsup _{r \rightarrow 0} \frac{\lambda(E \cap B(x, r))}{\lambda(B(x, r))}  \tag{7.17}\\
& \underline{D}(E, x)=\liminf _{r \rightarrow 0} \frac{\lambda(E \cap B(x, r))}{\lambda(B(x, r))} . \tag{7.18}
\end{align*}
$$

7.13. Theorem (Lebesgue Density Theorem). If $E \subset \mathbb{R}^{n}$ is a Lebesgue measurable set, then

$$
D(E, x)=1 \quad \text { for } \lambda \text {-almost all } x \in E \text {, }
$$

and

$$
D(E, x)=0 \quad \text { for } \lambda \text {-almost all } x \in \widetilde{E} .
$$

Proof. For the first part, let $B(r)$ denote the open ball centered at the origin of radius $r$ and let $f=\chi_{E \cap B(r)}$. Since $E \cap B(r)$ is bounded, it follows that $f$ is integrable and then Theorem 7.11 implies that $D(E \cap B(r), x)=1$ for $\lambda$-almost all $x \in E \cap B(r)$. Since $r$ is arbitrary, the result follows.

For the second part, take $f=\chi_{\widetilde{E} \cap B(r)}$ and conclude, as above, that $D(\widetilde{E} \cap$ $B(r), x)=1$ for $\lambda$-almost all $x \in \widetilde{E} \cap B(r)$. Observe that $D(E, x)=0$ for all such $x$, and thus the result follows since $r$ is arbitrary.

## Exercises for Section 7.2

1. Let $f$ be a measurable function defined on $\mathbb{R}^{n}$ with the property that, for some constant $C, f(x) \geq C|x|^{-n}$ for $|x| \geq 1$. Prove that $f$ is not integrable on $\mathbb{R}^{n}$. Hint: One way to proceed is to use Theorem 6.59.
2. Let $f \in L^{1}\left(\mathbb{R}^{n}\right)$ be a function that does not vanish identically on $B(0,1)$. Show that $M f \notin L^{1}\left(\mathbb{R}^{n}\right)$ by establishing the following inequality for all $x$ with $|x|>1$ :

$$
\begin{aligned}
M f(x) \geq f_{B(x,|x|+1)}|f| d \lambda & \geq \frac{1}{C(|x|+1)^{n}} \int_{B(0,1)}|f| d \lambda \\
& >\frac{1}{2^{n} C|x|^{n}} \int_{B(0,1)}|f| d \lambda
\end{aligned}
$$

where $C=\lambda[B(0,1)]$.
3. Prove that the maximal function $M f$ is not integrable on $\mathbb{R}^{n}$ unless $f$ is identically 0 (cf. previous Exercise 2).
4. Let $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$. Prove for each fixed $r>0$, that

$$
\int_{B(x, r)}|f| d \lambda
$$

is a continuous function of $x$.

### 7.3. The Radon-Nikodym Derivative - Another View

> We return to the concept of the Radon-Nikodym derivative in the setting of Lebesgue measure on $\mathbb{R}^{n}$. In this section it is shown that the RadonNikodym. derivative can be interpreted as a classical limiting process, very similar to that of the derivative of a function.

We now turn to the question of relating the derivative in the sense of (7.4) to the Radon-Nikodym derivative. Consider a $\sigma$-finite measure $\mu$ on $\mathbb{R}^{n}$ that is absolutely continuous with respect to Lebesgue measure. The Radon-Nikodym Theorem asserts the existence of a measurable function $f$ (the Radon-Nikodym derivative) such that $\mu$ can be represented as

$$
\mu(E)=\int_{E} f(y) d \lambda(y)
$$

for every Lebesgue measurable set $E \subset \mathbb{R}^{n}$. Theorem 7.11 implies that

$$
\begin{equation*}
D_{\lambda} \mu(x)=\lim _{r \rightarrow 0} \frac{\mu[B(x, r)]}{\lambda[B(x, r)]}=f(x) \tag{7.19}
\end{equation*}
$$

for $\lambda$-a.e. $x \in \mathbb{R}^{n}$. Thus, the Radon-Nikodym derivative of $\mu$ with respect to $\lambda$ and $D_{\lambda} \mu$ agree almost everywhere. Now we turn to measures that are singular with respect to Lebesgue measure.
7.14. Theorem. Let $\sigma$ be a Radon measure that is singular with respect to $\lambda$. Then

$$
D_{\lambda} \sigma(x)=0
$$

for $\lambda$-almost all $x \in \mathbb{R}^{n}$.
Proof. Since $\sigma \perp \lambda$ we know that $\sigma$ is concentrated on a Borel set $A$ with $\sigma(\widetilde{A})=\lambda(A)=0$. For each positive integer $k$, let

$$
E_{k}=\widetilde{A} \cap\left\{x: \limsup _{r \rightarrow 0} \frac{\sigma[B(x, r)]}{\lambda[B(x, r)]}>\frac{1}{k}\right\}
$$

In view of Exercise 4.8, we see for fixed $r$, that $\sigma[B(x, r)]$ is lower semicontinuous and therefore that $E_{k}$ is a Borel set. It suffices to show that

$$
\lambda\left(E_{k}\right)=0 \quad \text { for all } \quad k
$$

because $D_{\lambda} \sigma(x)=0$ for all $x \in \tilde{A}-\cup_{k=1}^{\infty} E_{k}$ and $\lambda(A)=0$. Referring to Theorems 4.63 and 4.52 , it follows that for every $\varepsilon>0$ there exists an open set $U_{\varepsilon} \supset E_{k}$ such
that $\sigma\left(U_{\varepsilon}\right)<\varepsilon$. For each $x \in E_{k}$ there exists a ball $B(x, r)$ with $0<r<1$ such that $B(x, r) \subset U_{\varepsilon}$ and $\lambda[B(x, r)]<k \sigma[B(x, r)]$. The collection of all such balls $B(x, r)$ provides a covering of $E_{k}$. Now employ Theorem 7.3 with $R=1$ to obtain a disjoint collection of balls, $\mathcal{F}$, such that

$$
E_{k} \subset \bigcup_{B \in \mathcal{F}} \widehat{B}
$$

Then

$$
\begin{aligned}
\lambda\left(E_{k}\right) & \leq \lambda\left\{\bigcup_{B \in \mathcal{F}} \widehat{B}\right\} \\
& \leq 5^{n} \sum_{B \in \mathcal{F}} \lambda(B) \\
& <5^{n} k \sum_{B \in \mathcal{F}} \sigma(B) \\
& \leq 5^{n} k \sigma\left(U_{\varepsilon}\right) \\
& \leq 5^{n} k \varepsilon .
\end{aligned}
$$

Since $\varepsilon$ is arbitrary, this shows that $\lambda\left(E_{k}\right)=0$.
This result together with (7.19) establishes the following theorem.
7.15. Theorem. Suppose $\nu$ is a Radon measure on $\mathbb{R}^{n}$. Let $\nu=\mu+\sigma$ be its Lebesgue decomposition with $\mu \ll \lambda$ and $\sigma \perp \lambda$. Finally, let $f$ denote the RadonNikodym derivative of $\mu$ with respect to $\lambda$. Then

$$
\lim _{r \rightarrow 0} \frac{\nu[B(x, r)]}{\lambda[B(x, r)]}=f(x)
$$

for $\lambda$-a.e. $x \in \mathbb{R}^{n}$.
7.16. Definition. Now we address the issue raised in Remark 7.2 concerning the use of concentric balls in the definition of (7.4). It can easily be shown that nonconcentric balls or even a more general class of sets could be used. For $x \in \mathbb{R}^{n}$, a sequence of Borel sets $\left\{E_{k}(x)\right\}$ is called a regular differentiation basis at $x$ provided there is a number $\alpha_{x}>0$ with the following property: There is a sequence of balls $B\left(x, r_{k}\right)$ with $r_{k} \rightarrow 0$ such that $E_{k}(x) \subset B\left(x, r_{k}\right)$ and

$$
\lambda\left(E_{k}(x)\right) \geq \alpha_{x} \lambda\left[B\left(x, r_{k}\right)\right] .
$$

The sets $E_{k}(x)$ are in no way related to $x$ except for the condition $E_{k} \subset B\left(x, r_{k}\right)$. In particular, the sets are not required to contain $x$.

The next result shows that Theorem 7.15 can be generalized to include regular differentiation bases.
7.17. Theorem. Suppose the hypotheses and notation of Theorem 7.15 are in force. Then for $\lambda$ almost every $x \in \mathbb{R}^{n}$, we have

$$
\lim _{k \rightarrow \infty} \frac{\sigma\left[E_{k}(x)\right]}{\lambda\left[E_{k}(x)\right]}=0
$$

and

$$
\lim _{k \rightarrow \infty} \frac{\mu\left[E_{k}(x)\right]}{\lambda\left[E_{k}(x)\right]}=f(x)
$$

whenever $\left\{E_{k}(x)\right\}$ is a regular differentiation basis at $x$.
Proof. In view of the inequalities

$$
\begin{equation*}
\frac{\alpha_{x} \sigma\left[E_{k}(x)\right]}{\lambda\left[E_{k}(x)\right]} \leq \frac{\sigma\left[E_{k}(x)\right]}{\lambda\left[B\left(x, r_{k}\right)\right]} \leq \frac{\sigma\left[B\left(x, r_{k}\right)\right]}{\lambda\left[B\left(x, r_{k}\right)\right]} \tag{7.20}
\end{equation*}
$$

the first conclusion of the Theorem follows from Lemma 7.14.
Concerning the second conclusion, Theorem 7.12 implies

$$
\lim _{r_{k} \rightarrow 0} f_{B\left(x, r_{k}\right)}|f(y)-f(x)| d \lambda(y)=0
$$

for almost all $x$ and consequently, by the same reasoning as in (7.20),

$$
\lim _{k \rightarrow \infty} f_{E_{k}(x)}|f(y)-f(x)| d \lambda(y)=0
$$

for almost all $x$. Hence, for almost all $x$ it follows that

$$
\lim _{k \rightarrow \infty} \frac{\mu\left[E_{k}(x)\right]}{\lambda\left[E_{k}(x)\right]}=\lim _{k \rightarrow \infty} f_{E_{k}(x)} f(y) d \lambda(y)=f(x)
$$

This leads immediately to the following theorem which is fundamental to the theory of functions of a single variable. For the companion result for functions of several variables, see Theorem 11.1. Also, see Exercise 2, Section 7.3, for a completely different proof.
7.18. Theorem. Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a nondecreasing function. Then $f^{\prime}(x)$ exists at $\lambda$-a.e. $x \in \mathbb{R}$.

Proof. Since $f$ is nondecreasing, Theorem 3.61 implies that $f$ is countinuous except possibly on the countable set $D=\left\{x_{1}, x_{2}, \ldots\right\}$. Indeed, from Theorem 3.61 we have

$$
f\left(x_{i}-\right)<f\left(x_{i}+\right)
$$

for each $x_{i} \in D$. Define $g: \mathbb{R} \rightarrow \mathbb{R}$ as $g(x)=f(x)$ if $x \notin D$ and $g(x)=f\left(x_{i}+\right)$ if $x_{i} \in D$. Then $g$ is a right continuous, nondecreasing function that agrees with $f$ except on the countable set $D$ (see exercise 1 Section 7.3). Now refer to Theorem 4.31 and theorem 4.32 to obtain a Borel measure $\mu$ such that

$$
\begin{equation*}
\mu((a, b])=g(b)-g(a) \tag{7.21}
\end{equation*}
$$

whenever $a<b$. For $x \in \mathbb{R}$ take as a regular differentiation basis an arbitrary sequence of half-open intervals $\left\{I_{k}(x)\right\}$ with $I_{k}=\left(x, x+h_{k}\right], h_{k}>0, h_{k} \rightarrow 0$. Indeed note that $\lambda\left(I_{k}(x)\right)=\frac{1}{2} \lambda\left[\left(x-h_{k}, x+h_{k}\right)\right]$. From Theorem 7.15 we have the decomposition $\mu=\tilde{\mu}+\sigma$ where $\tilde{\mu} \ll \lambda$ and $\sigma \ll \lambda$. Theorem 7.17 and Theorem 7.14 state that $D_{\lambda} \mu(x)=D_{\lambda} \tilde{\mu}(x)(x)+D_{\lambda} \sigma(x)=D_{\lambda} \tilde{\mu}(x)(x)$, for $\lambda$-almost every $x$. Hence, for $\lambda$-a.e. $x$, there exists $c_{x}$ such that

$$
\lim _{k \rightarrow \infty} \frac{\mu\left(I_{k}(x)\right)}{\lambda\left(I_{k}(x)\right)}=\lim _{k \rightarrow \infty} \frac{\mu\left(\left(x, x+h_{k}\right]\right)}{\lambda\left(\left(x, x+h_{k}\right]\right)}=c_{x}
$$

Clearly, the limit above is true for any sequence $h_{k} \rightarrow 0, h_{k}>0$, and thus

$$
\begin{equation*}
\lim _{h \rightarrow 0^{+}} \frac{\mu((x, x+h])}{\lambda((x, x+h])}=c_{x} \text { for } \lambda \text {-a.e. } x \tag{7.22}
\end{equation*}
$$

In a similar way we see that

$$
\begin{equation*}
\lim _{h \rightarrow 0^{+}} \frac{\mu((x-h, x])}{\lambda((x-h, x])}=c_{x}, \text { for } \lambda \text {-a.e. } x . \tag{7.23}
\end{equation*}
$$

From (7.22), (7.23) and (7.21) we have, for $\lambda$-a.e. $x$

$$
\begin{equation*}
\lim _{h \rightarrow 0^{+}} \frac{g(x+h)-g(x)}{h}=c_{x}=\lim _{h \rightarrow 0^{+}} \frac{g(x)-g(x-h)}{h} \tag{7.24}
\end{equation*}
$$

which means that $g^{\prime}(x)$ exists for $\lambda$-a.e. $x$ and $g^{\prime}(x)=c_{x}$ for such $x$. Let $G:=$ $\left\{x \in \mathbb{R}: g^{\prime}(x)\right.$ exists and $g(x)=f(x)$. Clearly, $\mathbb{R} \backslash G$ is a set of $\lambda$-measure zero. We now show that $f$ is differentiable at each $x \in G$ and $f^{\prime}(x)=g^{\prime}(x)$. Consider the sequence $h_{k} \rightarrow 0, h_{k}>0$. For each $h_{k}$ choose $0<h_{k}^{\prime}<h_{k}<h_{k}^{\prime \prime}$ such that $\frac{h_{k}^{\prime \prime}}{h_{k}} \rightarrow 1$ and $\frac{h_{k}^{\prime}}{h_{k}}=1$. Then

$$
\frac{h_{k}^{\prime}}{h_{k}} \cdot \frac{f\left(x+h_{k}^{\prime}\right)-f(x)}{h_{k}^{\prime}} \leq \frac{f\left(x+h_{k}\right)-f(x)}{h_{k}} \leq \frac{f\left(x+h_{k}^{\prime \prime}\right)-f(x)}{h_{k}^{\prime \prime}} \cdot \frac{h_{k}^{\prime \prime}}{h_{k}}
$$

Note that $h_{k}^{\prime}$ and $h_{k}^{\prime \prime}$ can be chosen so that $f\left(x+h_{k}^{\prime}\right)=g\left(x+h_{k}^{\prime}\right)$ and $f\left(x+h_{k}^{\prime \prime}\right)=$ $g\left(x+h_{k}^{\prime \prime}\right)$. Since $g(x)=f(x)$ for $\in G$, we obtain

$$
\frac{h_{k}^{\prime}}{h_{k}} \cdot \frac{g\left(x+h_{k}^{\prime}\right)-g(x)}{h_{k}^{\prime}} \leq \frac{f\left(x+h_{k}\right)-f(x)}{h_{k}} \leq \frac{g\left(x+h_{k}^{\prime \prime}\right)-g(x)}{h_{k}^{\prime \prime}} \cdot \frac{h_{k}^{\prime \prime}}{h_{k}} .
$$

Letting $h_{k} \rightarrow 0$ we obtain

$$
\begin{equation*}
\lim _{h_{k} \rightarrow 0^{+}} \frac{f\left(x+h_{k}\right)-f(x)}{h_{k}}=g^{\prime}(x) \tag{7.25}
\end{equation*}
$$

The same argument shows that

$$
\begin{equation*}
\lim _{h_{k} \rightarrow 0^{+}} \frac{f(x)-f\left(x-h_{k}\right)}{h_{k}}=g^{\prime}(x) \tag{7.26}
\end{equation*}
$$

Since (7.25) and (7.26) hold for any $h_{k} \rightarrow 0, h_{k}>0$ we conclude that $f^{\prime}(x)=$ $g^{\prime}(x)=c_{x}$ for each $x \in G$.

Another consequence of the above results is the following theorem concerning the derivative of the indefinite integral.
7.19. Theorem. Suppose $f$ is a Lebesgue integrable function defined on $[a, b]$. For each $x \in[a, b]$ let

$$
F(x)=\int_{a}^{x} f(t) d \lambda(t) .
$$

Then $F^{\prime}=f$ almost everywhere on $[a, b]$.
Proof. The derivative $F^{\prime}(x)$ is given by

$$
F^{\prime}(x)=\lim _{h \rightarrow 0} \frac{1}{h} \int_{x}^{x+h} f(t) d \lambda(t)
$$

Let $\mu$ be the measure defined by

$$
\mu(E)=\int_{E} f d \lambda
$$

for every measurable set $E$. Using intervals of the form $I_{h}(x)=[x, x+h]$ as a regular differentiation basis, it follows from Theorem 7.17 that

$$
\lim _{h \rightarrow 0} \frac{1}{h} \int_{x}^{x+h} f(t) d \lambda(t)=\lim _{h \rightarrow 0} \frac{\mu\left[I_{h}(x)\right]}{\lambda\left[I_{h}(x)\right]}=f(x)
$$

for almost all $x \in[a, b]$.

## Exercises for Section 7.3

1. Show that $g: \mathbb{R} \rightarrow \mathbb{R}$ defined in Theorem 7.18 is a nondecreasing and rightcontinuous function.
2. Here is an outline of an alternative proof of Theorem 7.18, which states that a nondecreasing function $f$ defined on $(a, b)$ is differentiable (Lebesgue) almost everywhere. For this, we introduce the Dini Derivatives:

$$
\begin{aligned}
D^{+} f\left(x_{0}\right) & =\limsup _{h \rightarrow 0^{+}} \frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h} \\
D_{+} f\left(x_{0}\right) & =\liminf _{h \rightarrow 0^{+}} \frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h} \\
D^{-} f\left(x_{0}\right) & =\limsup _{h \rightarrow 0^{-}} \frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h} \\
D_{-} f\left(x_{0}\right) & =\liminf _{h \rightarrow 0^{-}} \frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h}
\end{aligned}
$$

If all four Dini derivatives are finite and equal, then $f^{\prime}\left(x_{0}\right)$ exists and is equal to the common value. Clearly, $D_{+} f\left(x_{0}\right) \leq D^{+} f\left(x_{0}\right)$ and $D_{-} f\left(x_{0}\right) \leq D^{-} f\left(x_{0}\right)$. To prove that $f^{\prime}$ exists almost everywhere, it suffices to show that the set

$$
\left\{x: D^{+} f(x)>D_{-} f(x)\right\}
$$

has measure zero. A similar argument would apply to any two Dini derivatives. For any two rationals $r, s>0$, let

$$
E_{r, s}=\left\{x: D^{+} f(x)>r>s>D_{-} f(x)\right\} .
$$

The proof reduces to showing that $E_{r, s}$ has measure zero. If $x \in E_{r, s}$, then there exist arbitrarily small positive $h$ such that

$$
\frac{f(x-h)-f(x)}{-h}<s
$$

Fix $\varepsilon>0$ and use the Vitali Covering Theorem to find a countable family of closed, disjoint intervals $\left[x_{k}-h_{k}, x_{k}\right], k=1,2, \ldots$, such that

$$
\begin{gathered}
f\left(x_{k}\right)-f\left(x_{k}-h_{k}\right)<s h_{k} \\
\lambda\left(E_{r, s} \bigcap\left(\bigcup_{k=1}^{\infty}\left[x_{k}-h_{k}, x_{k}\right]\right)=\lambda\left(E_{r, s}\right),\right. \\
\sum_{k=1}^{m} h_{k}<(1+\varepsilon) \lambda\left(E_{r, s}\right)
\end{gathered}
$$

From this it follows that

$$
\sum_{k=1}^{\infty}\left[f\left(x_{k}\right)-f\left(x_{k}-h_{k}\right)\right]<s(1+\varepsilon) \lambda\left(E_{r, s}\right)
$$

For each point

$$
y \in A:=E_{r, s} \bigcap\left(\bigcup_{k=1}^{\infty}\left(x_{k}-h_{k}, x_{k}\right)\right),
$$

there exist arbitrarily small $h>0$ such that

$$
\frac{f(y+h)-f(y)}{h}>r .
$$

Employ the Vitali Covering Theorem again to obtain a countable family of disjoint, closed intervals $\left[y_{j}, y_{j}+h_{j}\right]$ such that

$$
\begin{aligned}
& \text { each }\left[y_{j}, y_{j}+h_{j}\right] \text { lies in some }\left[x_{k}-h_{k}, x_{k}\right], \\
& \qquad f\left(y_{j}+h_{j}\right)-f\left(y_{j}\right)>r h_{j} j=1,2, \ldots \\
& \sum_{j=1}^{\infty} h_{j} \geq \lambda(A)
\end{aligned}
$$

Since $f$ is nondecreasing, it follows that

$$
\sum_{j=1}^{\infty}\left[f\left(y_{j}+h_{j}\right)-f\left(y_{j}\right)\right] \leq \sum_{k=1}^{\infty}\left[f\left(x_{k}\right)-f\left(x_{k}-h_{k}\right)\right]
$$

and from this a contradiction is readily reached.

### 7.4. Functions of Bounded Variation

The main objective of this and the next section is to completely determine the conditions under which the following equation holds on an interval $[a, b]$ :

$$
f(x)-f(a)=\int_{a}^{x} f^{\prime}(t) d t \quad \text { for } \quad a \leq x \leq b
$$

This formula is well known in the context of Riemann integration and our purpose is to investigate its validity via the Lebesgue integral. It will be shown that the formula is valid precisely for the class of absolutely continuous functions. In this section we begin by introducing functions of bounded variation.

In the elementary version of the Fundamental Theorem of Calculus, it is assumed that $f^{\prime}$ exists at every point of $[a, b]$ and that $f^{\prime}$ is continuous. Since the Lebesgue integral is more general than the Riemann integral, one would expect a more general version of the Fundamental Theorem in the Lebesgue theory. What then would be the necessary assumptions? Perhaps it would be sufficient to assume that $f^{\prime}$ exists almost everywhere on $[a, b]$ and that $f^{\prime} \in L^{1}$. But this is obviously not true in view of the Cantor-Lebesgue function, $f$; see Example 5.7. We have seen that it is continuous, nondecreasing on $[0,1]$, and constant on each interval in the complement of the Cantor set. Consequently, $f^{\prime}=0$ at each point of the complement and thus

$$
1=f(1)-f(0)>\int_{0}^{1} f^{\prime}(t) d t=0
$$

The quantity $f(1)-f(0)$ indicates how much the function varies on $[0,1]$. Intuitively, one might have guessed that the quantity

$$
\int_{0}^{1}\left|f^{\prime}\right| d \lambda
$$

provides a measurement of the variation of $f$. Although this is false in general, for what class of functions is it true? We will begin to investigate the ideas surrounding these questions by introducing functions of bounded variation.
7.20. Definitions. Suppose a function $f$ is defined on $I=[a, b]$. The total variation of $f$ from $a$ to $x, x \leq b$, is defined by

$$
V_{f}(a ; x)=\sup \sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|
$$

where the supremum is taken over all finite sequences $a=t_{0}<t_{1}<\cdots<t_{k}=x$. $f$ is said to be of bounded variation (abbreviated, $B V$ ) on $[a, b]$ if $V_{f}(a ; b)<\infty$. If there is no danger of confusion, we will sometimes write $V_{f}(x)$ in place of $V_{f}(a ; x)$.

Note that if $f$ is of bounded variation on $[a, b]$ and $x \in[a, b]$, then

$$
|f(x)-f(a)| \leq V_{f}(a ; x) \leq V_{f}(a ; b)
$$

from which we see that $f$ is bounded.

It is easy to see that a bounded function that is either nonincreasing or nondecreasing is of bounded variation. Also, the sum (or difference) of two functions of bounded variation is again of bounded variation. The converse, which is not so immediate, is also true.
7.21. Theorem. Suppose $f$ is of bounded variation on $[a, b]$. Then $f$ can be written as

$$
f=f_{1}-f_{2}
$$

where both $f_{1}$ and $f_{2}$ are nondecreasing.

Proof. Let $x_{1}<x_{2} \leq b$ and let $a=t_{0}<t_{1}<\cdots<t_{k}=x_{1}$. Then

$$
\begin{equation*}
V_{f}\left(x_{2}\right) \geq\left|f\left(x_{2}\right)-f\left(x_{1}\right)\right|+\sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right| \tag{7.27}
\end{equation*}
$$

Now,

$$
V_{f}\left(x_{1}\right)=\sup \sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|
$$

over all sequences $a=t_{0}<t_{1}<\cdots<t_{k}=x_{1}$. Hence,

$$
\begin{equation*}
V_{f}\left(x_{2}\right) \geq\left|f\left(x_{2}\right)-f\left(x_{1}\right)\right|+V_{f}\left(x_{1}\right) \tag{7.28}
\end{equation*}
$$

In particular,

$$
V_{f}\left(x_{2}\right)-f\left(x_{2}\right) \geq V_{f}\left(x_{1}\right)-f\left(x_{1}\right) \quad \text { and } \quad V_{f}\left(x_{2}\right)+f\left(x_{2}\right) \geq V_{f}\left(x_{1}\right)+f\left(x_{1}\right)
$$

This shows that $V_{f}-f$ and $V_{f}+f$ are nondecreasing functions. The assertions thus follow by taking

$$
f_{1}=\frac{1}{2}\left(V_{f}+f\right) \quad \text { and } \quad f_{2}=\frac{1}{2}\left(V_{f}-f\right)
$$

7.22. Theorem. Suppose $f$ is of bounded variation on $[a, b]$. Then $f$ is Borel measurable and has at most a countable number of discontinuities. Furthermore, $f^{\prime}$ exists almost everywhere on $[a, b], f^{\prime}$ is Lebesgue measurable,

$$
\begin{equation*}
\left|f^{\prime}(x)\right|=V^{\prime}(x) \tag{7.29}
\end{equation*}
$$

for a.e. $x \in[a, b]$, and

$$
\begin{equation*}
\int_{a}^{b}\left|f^{\prime}(x)\right| d \lambda(t) \leq V_{f}(b) \tag{7.30}
\end{equation*}
$$

In particular, if $f$ is nondecreasing on $[a, b]$, then

$$
\begin{equation*}
\int_{a}^{b} f^{\prime}(x) d \lambda(x) \leq f(b)-f(a) \tag{7.31}
\end{equation*}
$$

Proof. We will first prove (7.31). Assume $f$ is nondecreasing and extend $f$ by defining $f(x)=f(b)$ for $x>b$ and for each positive integer $i$, let $g_{i}$ be defined by

$$
g_{i}(x)=i[f(x+1 / i)-f(x)] .
$$

Since $f$ in nondecreasing it follows that $g_{i}$ is a Borel function. Consequently, the functions $u$ and $v$ defined by

$$
\begin{align*}
u(x) & =\limsup _{i \rightarrow \infty} g_{i}(x)  \tag{7.32}\\
v(x) & =\liminf _{i \rightarrow \infty} g_{i}(x)
\end{align*}
$$

are also Borel functions. We know from Theorem 7.18 that $f^{\prime}$ exists a.e. Hence, it follows that $f^{\prime}=u$ a.e. and is therefore Lebesgue measurable. Now, each $g_{i}$ is nonnegative because $f$ is nondecreasing and therefore we may employ Fatou's lemma to conclude

$$
\begin{aligned}
\int_{a}^{b} f^{\prime}(x) d \lambda(x) & \leq \liminf _{i \rightarrow \infty} \int_{a}^{b} g_{i}(x) d \lambda(x) \\
& =\liminf _{i \rightarrow \infty} i \int_{a}^{b}[f(x+1 / i)-f(x)] d \lambda(x) \\
& =\liminf _{i \rightarrow \infty} i\left[\int_{a+1 / i}^{b+1 / i} f(x) d \lambda(x)-\int_{a}^{b} f(x) d \lambda(x)\right] \\
& =\liminf _{i \rightarrow \infty} i\left[\int_{b}^{b+1 / i} f(x) d \lambda(x)-\int_{a}^{a+1 / i} f(x) d \lambda(x)\right] \\
& \leq \liminf _{i \rightarrow \infty} i\left[\frac{f(b+1 / i)}{i}-\frac{f(a)}{i}\right] \\
& =\liminf _{i \rightarrow \infty} i\left[\frac{f(b)}{i}-\frac{f(a)}{i}\right]=f(b)-f(a)
\end{aligned}
$$

In establishing the last inequality, we have used the fact that $f$ is nondecreasing.
Now suppose that $f$ is an arbitrary function of bounded variation. Since $f$ can be written as the difference of two nondecreasing functions, Theorem 7.21, it follows from Theorem 3.61, that the set $D$ of discontinuities of $f$ is countable. For each real number $t$ let $A_{t}:=\{f>t\}$. Then

$$
(a, b) \cap A_{t}=\left((a, b) \cap\left(A_{t}-D\right)\right) \cup\left((a, b) \cap A_{t} \cap D\right) .
$$

The first set on the right is open since $f$ is continuous at each point of $(a, b)-D$. Since $D$ is countable, the second set is a Borel set; therefore, so is $(a, b) \cap A_{t}$ which implies that $f$ is a Borel function.

The statements in the Theorem referring to the almost everywhere differentiability of $f$ follows from Theorem 7.18; the measurability of $f^{\prime}$ is addressed in (7.32).

Similarly, since $V_{f}$ is a nondecreasing function, we have that $V_{f}^{\prime}$ exists almost everywhere. Furthermore, with $f=f_{1}-f_{2}$ as in the previous theorem and recalling that $f_{1}^{\prime}, f_{2}^{\prime} \geq 0$ almost everywhere, it follows that

$$
\left|f^{\prime}\right|=\left|f_{1}^{\prime}-f_{2}^{\prime}\right| \leq\left|f_{1}^{\prime}\right|+\left|f_{2}^{\prime}\right|=f_{1}^{\prime}+f_{2}^{\prime}=V_{f}^{\prime} \quad \text { almost everywhere on }[a, b] .
$$

To prove (7.29) we will show that

$$
E:=[a, b] \bigcap\left\{t: V_{f}^{\prime}(t)>\left|f^{\prime}(t)\right|\right\}
$$

has measure zero. For each positive integer $m$ let $E_{m}$ be the set of all $t \in E$ such that $\tau_{1} \leq t \leq \tau_{2}$ with $0<\tau_{2}-\tau_{1}<\frac{1}{m}$ implies

$$
\begin{equation*}
\frac{V_{f}\left(\tau_{2}\right)-V_{f}\left(\tau_{1}\right)}{\tau_{2}-\tau_{1}}>\frac{\left|f\left(\tau_{2}\right)-f\left(\tau_{1}\right)\right|}{\tau_{2}-\tau_{1}}+\frac{1}{m} . \tag{7.33}
\end{equation*}
$$

Since each $t \in E$ belongs to $E_{m}$ for sufficiently large $m$ we see that

$$
E=\bigcup_{m=1}^{\infty} E_{m}
$$

and thus it suffices to show that $\lambda\left(E_{m}\right)=0$ for each $m$. Fix $\varepsilon>0$ and let $a=t_{0}<t_{1}<\cdots<t_{k}=b$ be a partition of $[a, b]$ such that $\left|t_{i}-t_{i-1}\right|<\frac{1}{m}$ for each $i$ and

$$
\begin{equation*}
\sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|>V_{f}(b)-\frac{\varepsilon}{m} \tag{7.34}
\end{equation*}
$$

For each interval in the partition, (7.28) states that

$$
\begin{equation*}
V_{f}\left(t_{i}\right)-V_{f}\left(t_{i-1}\right) \geq\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right| \tag{7.35}
\end{equation*}
$$

while (7.33) implies

$$
\begin{equation*}
V_{f}\left(t_{i}\right)-V_{f}\left(t_{i-1}\right) \geq\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|+\frac{t_{i}-t_{i-1}}{m} \tag{7.36}
\end{equation*}
$$

if the interval contains a point of $E_{m}$. Let $\mathcal{F}_{1}$ denote those intervals of the partition that do not contain any points of $E_{m}$ and let $\mathcal{F}_{2}$ denote those intervals that do contain points of $E_{m}$. Then, since $\lambda\left(E_{m}\right) \leq \sum_{I \in \mathcal{F}_{2}} b_{I}-a_{I}$,

$$
\begin{align*}
V_{f}(b) & =\sum_{i=1}^{k} V_{f}\left(t_{i}\right)-V_{f}\left(t_{i-1}\right) \\
& =\sum_{I \in \mathcal{F}_{1}} V_{f}\left(b_{I}\right)-V_{f}\left(a_{I}\right)+\sum_{I \in \mathcal{F}_{2}} V_{f}\left(b_{I}\right)-V_{f}\left(a_{I}\right) \\
& =\sum_{I \in \mathcal{F}_{1}} f\left(b_{I}\right)-f\left(a_{I}\right)+\sum_{I \in \mathcal{F}_{2}} f\left(b_{I}\right)-f\left(a_{I}\right)+\frac{b_{I}-a_{I}}{m} \quad \text { by }(7.35) \text { and }(7.36) \\
& \geq \sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|+\frac{\lambda\left(E_{m}\right)}{m} \\
& \geq V_{f}(b)-\frac{\varepsilon}{m}+\frac{\lambda\left(E_{m}\right)}{m}, \quad \text { by }(7.34) \tag{7.34}
\end{align*}
$$

and therefore $\lambda\left(E_{m}\right) \leq \varepsilon$, from which we conclude that $\lambda\left(E_{m}\right)=0$ since $\varepsilon$ is arbitrary. Thus (7.29) is established.

Finally we apply (7.29) and (7.31) to obtain

$$
\int_{a}^{b}\left|f^{\prime}\right| d \lambda=\int_{a}^{b} V^{\prime} d \lambda \leq V_{f}(b)-V_{f}(a)=V_{f}(b)
$$

and the proof is complete.

## Exercises for Section 7.4

1. Prove that a function of bounded variation is a Borel measurable function.
2. If $f$ is of bounded variation on $[a, b]$, Theorem 7.21 states that $f=f_{1}-f_{2}$ where both $f_{1}$ and $f_{2}$ are nondecreasing. Prove that

$$
f_{1}(x)=\sup \left\{\sum_{i=1}^{k}\left(f\left(t_{i}\right)-f\left(t_{i-1}\right)\right)^{+}\right\}
$$

where the supremum is taken over all partitions $a=t_{0}<t_{1}<\cdots<t_{k}=x$.

### 7.5. The Fundamental Theorem of Calculus

We introduce absolutely continuous functions and show that they are precisely those functions for which the Fundamental Theorem of Calculus is valid.
7.23. Definition. A function $f$ defined on an interval $I=[a, b]$ is said to be absolutely continuous on $I$ (briefly, AC on $I$ ) if for every $\varepsilon>0$ there exists $\delta>0$ such that

$$
\sum_{i=1}^{k}\left|f\left(b_{i}\right)-f\left(a_{i}\right)\right|<\varepsilon
$$

for any finite collection of nonoverlapping intervals $\left[a_{1}, b_{1}\right],\left[a_{2}, b_{2}\right], \ldots,\left[a_{k}, b_{k}\right]$ in $I$ with

$$
\sum_{i=1}^{k}\left|b_{i}-a_{i}\right|<\delta
$$

Observe if $f$ is AC , then it is easy to show that

$$
\sum_{i=1}^{\infty}\left|f\left(b_{i}\right)-f\left(a_{i}\right)\right| \leq \varepsilon
$$

for any countable collection of nonoverlapping intervals with

$$
\begin{equation*}
\sum_{i=1}^{\infty}\left|b_{i}-a_{i}\right|<\delta \tag{7.37}
\end{equation*}
$$

Indeed, if $f$ is AC , then (7.37) holds for any partial sum, and therefore for the limit of the partial sums. Thus, it holds for the whole series.

From the definition it follows that an absolutely continuous function is uniformly continuous. The converse is not true as we shall see illustrated later by the CantorLebesgue function. (Of course it can be shown directly that the Cantor-Lebesgue function is not absolutely continuous: see Exercise 1, Section 7.5). The reader can easily verify that any Lipschitz function is absolutely continuous. Another example of an AC function is given by the indefinite integral; let $f$ be an integrable function on $[a, b]$ and set

$$
\begin{equation*}
F(x)=\int_{a}^{x} f(t) d \lambda(t) \tag{7.38}
\end{equation*}
$$

For any nonoverlapping collection of intervals in $[a, b]$ we have

$$
\sum_{i=1}^{k}\left|F\left(b_{i}\right)-F\left(a_{i}\right)\right| \leq \int_{\cup\left[a_{i}, b_{i}\right]}|f| d \lambda
$$

With the help of Exercise 1, Section 6.6, we know that the set function $\mu$ defined by

$$
\mu(E)=\int_{E}|f| d \lambda
$$

is a measure, and clearly it is absolutely continuous with respect to Lebesgue measure. Referring to Theorem 6.40, we see that $F$ is absolutely continuous.
7.24. Notation. The following notation will be used frequently throughout. If $I \subset \mathbb{R}$ is an interval, we will denote its endpoints by $a_{I}, b_{I}$; thus, if $I$ is closed, $I=\left[a_{I}, b_{I}\right]$.
7.25. Theorem. An absolutely continuous function on $[a, b]$ is of bounded variation.

Proof. Let $f$ be absolutely continuous. Choose $\varepsilon=1$ and let $\delta>0$ be the corresponding number provided by the definition of absolute continuity. Subdivide [a,b] into a finite collection $\mathcal{F}$ of nonoverlapping subintervals $I=\left[a_{I}, b_{I}\right]$ each of whose length is less than $\delta$. Then,

$$
\left|f\left(b_{I}\right)-f\left(a_{I}\right)\right|<1
$$

for each $I \in \mathcal{F}$. Consequently, if $\mathcal{F}$ consists of $M$ elements, we have

$$
\sum_{I \in \mathcal{F}}\left|f\left(b_{I}\right)-f\left(a_{I}\right)\right|<M
$$

To show that $f$ is of bounded variation on [a,b], consider an arbitrary partition $a=t_{0}<t_{1}<\cdots<t_{k}=b$. Since the sum

$$
\sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|
$$

is not decreased by adding more points to this partition, we may assume each interval of this partition is a subset of some $I \in \mathcal{F}$. But then, for each $I \in \mathcal{F}$, it follows that

$$
\sum_{i=1}^{k} \chi_{I}\left(t_{i}\right) \chi_{I}\left(t_{i-1}\right)\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|<1
$$

(this is simply saying that the sum is taken over only those intervals $\left[t_{i-1}, t_{i}\right]$ that are contained in $I$ ) and consequently,

$$
\sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|<M
$$

thus proving that the total variation of $f$ on $[\mathrm{a}, \mathrm{b}]$ is no more than $M$.

Next, we introduce a property that is of great importance concerning absolutely continuous functions. Later we will see that this property is one among three that characterize absolutely continuous functions (see Corollary 7.36). This concept is due to Lusin, now called "condition $N$."
7.26. Definition. A function $f$ defined on $[a, b]$ is said to satisfy condition $N$ if $f$ preserves sets of Lebesgue measure zero; that is, $\lambda[f(E)]=0$ whenever $E \subset[a, b]$ with $\lambda(E)=0$.
7.27. Theorem. If $f$ is an absolutely continuous function on $[a, b]$, then $f$ satisfies condition $N$.

Proof. Choose $\varepsilon>0$ and let $\delta>0$ be the corresponding number provided by the definition of absolute continuity. Let $E$ be a set of measure zero. Then there is an open set $U \supset E$ with $\lambda(U)<\delta$. Since $U$ is the union of a countable collection $\mathcal{F}$ of disjoint open intervals, we have

$$
\sum_{I \in \mathcal{F}} \lambda(I)<\delta
$$

The closure of each interval $I$ contains an interval $I^{\prime}=\left[a_{I^{\prime}}, b_{I^{\prime}}\right]$ at whose endpoints $f$ assumes its maximum and minimum on the closure of $I$. Then

$$
\lambda[f(I)]=\left|f\left(b_{I^{\prime}}\right)-f\left(a_{I^{\prime}}\right)\right|
$$

and the absolute continuity of $f$ along with (7.37) imply

$$
\lambda[f(E)] \leq \sum_{I \in \mathcal{F}} \lambda[f(I)]=\sum_{I \in \mathcal{F}}\left|f\left(b_{I^{\prime}}\right)-f\left(a_{I^{\prime}}\right)\right|<\varepsilon
$$

Since $\varepsilon$ is arbitrary, this shows that $f(E)$ has measure zero.
7.28. Remark. This result shows that the Cantor-Lebesgue function is not absolutely continuous, since it maps the Cantor set (of Lebesgue measure zero) onto $[0,1]$. Thus, there are continuous functions of bounded variation that are not absolutely continuous.
7.29. Theorem. Suppose $f$ is an arbitrary function defined on $[a, b]$. Let

$$
E_{f}:=(a, b) \cap\left\{x: f^{\prime}(x) \text { exists and } f^{\prime}(x)=0\right\} .
$$

Then $\lambda\left[f\left(E_{f}\right)\right]=0$.

## Proof. Step 1:

Initially, we will assume $f$ is bounded; let $M:=\sup \{f(x): x \in[a, b]\}$ and $m:=$ $\inf \{f(x): x \in[a, b]\}$. Choose $\varepsilon>0$. For each $x \in E_{f}$ there exists $\delta=\delta(x)>0$ such that

$$
|f(x+h)-f(x)|<\varepsilon h
$$

and

$$
|f(x-h)-f(x)|<\varepsilon h
$$

whenever $0<h<\delta(x)$. Thus, for each $x \in E_{f}$ we have a collection of intervals of the form $[x-h, x+h], 0<h<\delta(x)$, with the property that for arbitrary points $a^{\prime}, b^{\prime} \in I=[x-h, x+h]$, then

$$
\begin{align*}
\left|f\left(b^{\prime}\right)-f\left(a^{\prime}\right)\right| & \leq\left|f\left(b^{\prime}\right)-f(x)\right|+\left|f\left(a^{\prime}\right)-f(x)\right| \\
& <\varepsilon\left|b^{\prime}-x\right|+\varepsilon\left|a^{\prime}-x\right|  \tag{7.39}\\
& \leq \varepsilon \lambda(I)
\end{align*}
$$

We will adopt the following notation: For each interval $I$ let $I^{\prime}$ be an interval such that $I^{\prime} \supset I$ has the same center as $I$ but is 5 times as long. Using the definition of Lebesgue measure, we can find an open set $U \supset E_{f}$ with the property $\lambda(U)-\lambda\left(E_{f}\right)<\varepsilon$ and let $\mathcal{G}$ be the collection of intervals $I$ such that $I^{\prime} \subset U$ and $I^{\prime}$ satisfies (7.39). Then appeal to Theorem 7.3 to find a disjoint subfamily $\mathcal{F} \subset \mathcal{G}$ such that

$$
E_{f} \subset \bigcup_{I \in \mathcal{F}} I^{\prime}
$$

Since $f$ is bounded, each interval $I^{\prime}$ that is associated with $I \in \mathcal{F}$ contains points $a_{I^{\prime}}, b_{I^{\prime}}$ such that $f\left(b_{I^{\prime}}\right)>M_{I^{\prime}}-\varepsilon \lambda\left(I^{\prime}\right)$ and $f\left(a_{I^{\prime}}\right)<m_{I^{\prime}}+\varepsilon \lambda\left(I^{\prime}\right)$ where $\varepsilon>0$ is chosen as above, Thus,

$$
f\left(I^{\prime}\right) \subset\left[m_{I^{\prime}}, M_{I^{\prime}}\right] \subset\left[f\left(a_{I^{\prime}}\right)+\varepsilon \lambda\left(I^{\prime}\right), f\left(b_{I^{\prime}}\right)-\varepsilon \lambda\left(I^{\prime}\right)\right]
$$

and thus, by (7.39)

$$
\lambda\left[f\left(I^{\prime}\right)\right] \leq\left|M_{I^{\prime}}-m_{I^{\prime}}\right| \leq\left|f\left(b_{I^{\prime}}\right)-f\left(a_{I^{\prime}}\right)\right|+2 \varepsilon \lambda\left(I^{\prime}\right)<3 \varepsilon \lambda\left(I^{\prime}\right)
$$

Then, using the fact that $\mathcal{F}$ is a disjoint family,

$$
\begin{aligned}
\lambda\left(f\left(E_{f}\right)\right) \subset \lambda\left[f\left(\bigcup_{I \in \mathcal{F}} I^{\prime}\right)\right] & \leq \bigcup_{I \in \mathcal{F}} \lambda\left(f\left(I^{\prime}\right)\right)<\sum_{I \in \mathcal{F}} \varepsilon \lambda(I) \\
& =5 \sum_{I \in \mathcal{F}} \varepsilon \lambda(I) \\
& \leq 5 \varepsilon \lambda(U)<5 \varepsilon\left(\lambda\left(E_{f}\right)+\varepsilon\right)
\end{aligned}
$$

Since $\varepsilon>0$ is arbitrary, we conclude that $\lambda\left(f\left(E_{f}\right)\right)=0$ as desired.

## Step 2:

Now assume $f$ is an arbitrary function on $[a, b]$. Let $H: \mathbb{R} \rightarrow[0,1]$ be a smooth, strictly increasing function with $H^{\prime}>0$ on $\mathbb{R}$. Note that both $H$ and $H^{-1}$ are absolutely continuous. Define $g:=H \circ f$. Then $g$ is bounded and $g^{\prime}(x)=H^{\prime}(f(x)) f^{\prime}(x)$ holds whenever either $g$ or $f$ is differentiable at $x$. Then $f^{\prime}(x)=0$ iff $g^{\prime}(x)=0$ and therefore $E_{f}=E_{g}$. From Step 1, we know $\lambda\left[g\left(E_{g}\right)\right]=0=\lambda\left[g\left(E_{f}\right)\right]$. But $f(A)=H^{-1}[g(A)]$ for any $A \subset \mathbb{R}$, we have $f\left(E_{f}\right)=H^{-1}\left[g\left(E_{f}\right)\right]$ and therefore $\lambda\left[f\left(E_{f}\right)\right]=0$ since $H^{-1}$ preserves sets of measure zero.
7.30. Theorem. If $f$ is absolutely continuous on $[a, b]$ with the property that $f^{\prime}=0$ almost everywhere, then $f$ is constant.

Proof. Let

$$
E=(a, b) \cap\left\{x: f^{\prime}(x)=0\right\}
$$

so that $[a, b]=E \cup N$ where $N$ is of measure zero. Then $\lambda[f(E)]=\lambda[f(N)]=0$ by the previous result and Theorem 7.27. Thus, since $f([a, b])$ is an interval and also of measure zero, $f$ must be constant.

We now have reached the main objective.
7.31. Theorem (The Fundamental Theorem of Calculus). $f:[a, b] \rightarrow \mathbb{R}$ is absolutely continuous if and only if $f^{\prime}$ exists a.e. on $(a, b), f^{\prime}$ is integrable on $(a, b)$, and

$$
f(x)-f(a)=\int_{a}^{x} f^{\prime}(t) d \lambda(t) \quad \text { for } \quad x \in[a, b]
$$

Proof. The sufficiency follows from integration theory as discussed in (7.38).
As for necessity, recall that $f$ is $B V$ ( Theorem 7.25), and therefore by Theorem 7.22 that $f^{\prime}$ exists almost everywhere and is integrable. Hence, it is meaningful to define

$$
F(x)=\int_{a}^{x} f^{\prime}(t) d \lambda(t)
$$

Then $F$ is absolutely continuous (as in (7.38)). By Theorem 7.19 we have that $F^{\prime}=f^{\prime}$ almost everywhere on $[a, b]$. Thus, $F-f$ is an absolutely continuous function whose derivative is zero almost everywhere. Therefore, by Theorem 7.30, $F-f$ is constant on $[a, b]$ so that $[F(x)-f(x)]=[F(a)-f(a)]$ for all $x \in[a, b]$. Since $F(a)=0$, we have $F=f-f(a)$.
7.32. Corollary. If $f$ is an absolutely continuous function on $[a, b]$ then the total variation function $V_{f}$ of $f$ is also absolutely continuous on $[a, b]$ and

$$
V_{f}(x)=\int_{a}^{x}\left|f^{\prime}\right| d \lambda
$$

for each $x \in[a, b]$.
Proof. We know that $V_{f}$ is a bounded, nondecreasing function and thus of bounded variation. From Theorem 7.22 we know that

$$
\int_{a}^{x}\left|f^{\prime}\right| d \lambda=\int_{a}^{x} V_{f}^{\prime} d \lambda \leq V_{f}(x)
$$

for each $x \in[a, b]$.
Fix $x \in[a, b]$ and let $\varepsilon>0$. Choose a partition $\left\{a=t_{0}<t_{1}<\cdots<t_{k}=x\right\}$ so that

$$
V_{f}(x) \leq \sum_{i=1}^{k}\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|+\varepsilon .
$$

In view of the Theorem 7.31

$$
\left|f\left(t_{i}\right)-f\left(t_{i-1}\right)\right|=\left|\int_{t_{i-1}}^{t_{i}} f^{\prime} d \lambda\right| \leq \int_{t_{i-1}}^{t_{i}}\left|f^{\prime}\right| d \lambda
$$

for each $i$. Thus

$$
V_{f}(x) \leq \sum_{i=1}^{k} \int_{t_{i-1}}^{t_{i}}\left|f^{\prime}\right| d \lambda+\varepsilon \leq \int_{a}^{x}\left|f^{\prime}\right| d \lambda+\varepsilon
$$

Since $\varepsilon$ is arbitrary we conclude that

$$
V_{f}(x)=\int_{a}^{x}\left|f^{\prime}\right| d \lambda .
$$

for each $x \in[a, b]$.

## Exercises for Section 7.5

1. Prove directly from the definition that the Cantor-Lebesgue function is not absolutely continuous.
2. Prove that a Lipschitz function on $[a, b]$ is absolutely continuous.
3. Prove directly from definitions that a Lipschitz function satisfies condition $N$.
4. Suppose $f$ is a Lipschitz function defined on $\mathbb{R}$ having Lipschitz constant $C$. Prove that $\lambda[f(A)] \leq C \lambda(A)$ whenever $A \subset \mathbb{R}$ is Lebesgue measurable.
5. Let $\left\{f_{k}\right\}$ be a uniformly bounded sequence of absolutely continuous functions on $[0,1]$. Suppose that $f_{k} \rightarrow f$ in $L^{1}[0,1]$ and that $\left\{f_{k}^{\prime}\right\}$ is Cauchy in $L^{1}[0,1]$. Prove that $f=g$ almost everywhere, where $g$ is absolutely continuous on $[0,1]$.
6 . Let $f$ be a strictly increasing continuous function defined on $(0,1)$. Prove that $f^{\prime}>0$ almost everywhere if and only if $f^{-1}$ is absolutely continuous.
6. Prove that the sum and product of absolutely continuous functions are absolutely continuous.
7. Prove that the composition of $B V$ functions is not necessarily $B V$. (Hint: Consider the composition of $f(x)=\sqrt{x}$ and $g(x)=x^{2} \cos ^{2}\left(\frac{\pi}{2 x}\right), x \neq 0, g(0)=0$, both defined on $[0,1])$.
8. Find two absolutely continuous functions $f, g:[0,1] \rightarrow[0,1]$ such that their composition is not absolutely continuous. However, show that if $g:[a, b] \rightarrow[c, d]$ is absolutely continuous and $f:[c, d] \rightarrow \mathbb{R}$ is Lipschitz then $f \circ g$ is absolutely continuous.
9. Establish the integration by parts formula: If $f$ and $g$ are absolutely continuous functions defined on $[a, b]$, then

$$
\int_{a}^{b} f^{\prime} g d x=f(b) g(b)-f(a) g(a)-\int_{a}^{b} f g^{\prime} d x
$$

11. Give an example of a function $f:(0,1) \rightarrow \mathbb{R}$ that is differentiable everywhere but is not absolutely continuous. Compare this with Theorem 7.50.
12. Given a Lebesgue measurable set $E$, prove that $\{x: D(E, x)=1\}$ is a Borel set.

### 7.6. Variation of Continuous Functions

One possibility of determining the variation of a function is the following. Consider the graph of $f$ in the $(x, y)$-plane and for each $y$, let $N(y)$ denote the number of times the horizontal line passing through $(0, y)$ intersects the graph of $f$. It seems plausible that

$$
\int_{\mathbb{R}} N(y) d \lambda(y)
$$

should equal the variation of $f$ on $[a, b]$. In case $f$ is a continuous nondecreasing function, this is easily seen to be true. The next theorem provides the general result. First, we introduce some notation: If $f: \mathbb{R} \rightarrow \mathbb{R}$ and $E \subset \mathbb{R}$, then

$$
\begin{equation*}
N(f, E, y) \tag{7.40}
\end{equation*}
$$

denotes the (possibly infinite) number of points in the set $E \cap f^{-1}\{y\}$. Thus, $N(f, E, y)$ is the number of points in $E$ that are mapped onto $y$.
7.33. Theorem. Let $f$ be a continuous function defined on $[a, b]$. Then $N(f,[a, b], y)$ is a Borel measurable function (of y) and

$$
V_{f}(b)=\int_{R^{1}} N(f,[a, b], y) d \lambda(y) .
$$

Proof. For brevity throughout the proof, we will simply write $N(y)$ for $N(f,[a, b], y)$.

Let $m \leq N(y)$ be a nonnegative integer and let $x_{1}, x_{2}, \cdots, x_{m}$ be points that are mapped into $y$. Thus, $\left\{x_{1}, x_{2}, \ldots, x_{m}\right\} \subset f^{-1}\{y\}$. For each positive integer $i$, consider a partition, $\mathcal{P}_{i}=\left\{a=t_{0}<t_{1}<\ldots<t_{k}=b\right\}$ of $[a, b]$ such that the
length of each interval $I$ is less than $1 / i$. Choose $i$ so large that each interval of $\mathcal{P}_{i}$ contains at most one $x_{j}, j=1,2, \ldots, m$. Then

$$
m \leq \sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)
$$

Consequently,

$$
\begin{equation*}
m \leq \liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)\right\} \tag{7.41}
\end{equation*}
$$

Since $m$ is an arbitrary positive integer with $m \leq N(y)$, we obtain

$$
\begin{equation*}
N(y) \leq \liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)\right\} \tag{7.42}
\end{equation*}
$$

On the other hand, for any partition $\mathcal{P}_{i}$ we obviously have

$$
\begin{equation*}
N(y) \geq \sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y) \tag{7.43}
\end{equation*}
$$

provided that each point of $f^{-1}(y)$ is contained in the interior of some interval $I \in \mathcal{P}_{i}$. Thus (7.43) holds for all but finitely many $y$ and therefore

$$
N(y) \geq \limsup _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)\right\}
$$

holds for all but countably many $y$. Hence, with (7.42), we have

$$
\begin{equation*}
N(y)=\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)\right\} \quad \text { for all but countably many } y \tag{7.44}
\end{equation*}
$$

Since $\sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}$ is a Borel measurable function, it follows that $N$ is also Borel measurable. For any interval $I \in P_{i}$,

$$
\lambda[f(I)]=\int_{R^{1}} \chi_{f(I)}(y) d \lambda(y)
$$

and therefore by (7.43),

$$
\begin{aligned}
\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)] & =\sum_{I \in \mathcal{P}_{i}} \int_{R^{1}} \chi_{f(I)}(y) d \lambda(y) \\
& =\int_{R^{1}} \sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y) d \lambda(y) \\
& \leq \int_{R^{1}} N(y) d \lambda(y)
\end{aligned}
$$

which implies

$$
\begin{equation*}
\limsup _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)]\right\} \leq \int_{R^{1}} N(y) d \lambda(y) \tag{7.45}
\end{equation*}
$$

For the opposite inequality, observe that Fatou's lemma and (7.44) yield

$$
\begin{aligned}
\liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)]\right\} & =\liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \int_{R^{1}} \chi_{f(I)}(y) d \lambda(y)\right\} \\
& =\liminf _{i \rightarrow \infty}\left\{\int_{R^{1}} \sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y) d \lambda(y)\right\} \\
& \geq \int_{R^{1}}\left\{\liminf _{i \rightarrow \infty} \sum_{I \in \mathcal{P}_{i}} \chi_{f(I)}(y)\right\} d \lambda(y) \\
& =\int_{R^{1}} N(y) d \lambda(y) .
\end{aligned}
$$

Thus, we have

$$
\begin{equation*}
\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)]\right\}=\int_{R^{1}} N(y) d \lambda(y) \tag{7.46}
\end{equation*}
$$

We will conclude the proof by showing that the limit on the left-side is equal to $V_{f}(b)$. First, recall the notation introduced in Notation 7.24: If $I$ is an interval belonging to a partition $\mathcal{P}_{i}$, we will denote the endpoints of this interval by $a_{I}, b_{I}$. Thus, $I=\left[a_{I}, b_{I}\right]$. We now proceed with the proof by selecting a sequence of partitions $\mathcal{P}_{i}$ with the property that each subinterval $I$ in $\mathcal{P}_{i}$ has length less than $\frac{1}{i}$ and

$$
\lim _{i \rightarrow \infty} \sum_{I \in \mathcal{P}_{i}}\left|f\left(b_{i}\right)-f\left(a_{i}\right)\right|=V_{f}(b)
$$

Then,

$$
\begin{aligned}
V_{f}(b) & =\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \mid f\left(b_{I}\right)-f\left(a_{I} \mid\right)\right\} \\
& \leq \liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)]\right\}
\end{aligned}
$$

We now show that

$$
\begin{equation*}
\limsup _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)]\right\} \leq V_{f}(b) \tag{7.47}
\end{equation*}
$$

which will conclude the proof. For this, let $I^{\prime}=\left[a_{I^{\prime}}, b_{I^{\prime}}\right]$ be an interval contained in $I=\left[a_{i}, b_{i}\right]$ such that $f$ assumes its maximum and minimum on $I$ at the endpoints of $I^{\prime}$. Let $\mathcal{Q}_{i}$ denote the partition formed by the endpoints of $I \in \mathcal{P}_{i}$ along with
the endpoints of the intervals $I^{\prime}$. Then

$$
\begin{aligned}
\sum_{I \in \mathcal{P}_{i}} \lambda[f(I)] & =\sum_{I \in \mathcal{P}_{i}}\left|f\left(b_{I^{\prime}}\right)-f\left(a_{I^{\prime}}\right)\right| \\
& \leq \sum_{I \in \mathcal{Q}_{i}}\left|f\left(b_{i}\right)-f\left(a_{i}\right)\right| \\
& \leq V_{f}(b)
\end{aligned}
$$

thereby establishing (7.47).
7.34. Corollary. Suppose $f$ is a continuous function of bounded variation on $[a, b]$. Then the total variation function $V_{f}(\cdot)$ is continuous on $[a, b]$. In addition, if $f$ also satisfies condition $N$, then so does $V_{f}(\cdot)$.

Proof. Fix $x_{0} \in[a, b]$. By the previous result

$$
V_{f}\left(x_{0}\right)=\int_{\mathbb{R}} N\left(f,\left[a, x_{0}\right], y\right) d \lambda(y)
$$

For $a \leq x_{0}<x \leq b$,

$$
N(f,[a, x], y)-N\left(f,\left[a, x_{0}\right], y\right)=N\left(f,\left(x_{0}, x\right], y\right)
$$

for each $y$ such that $N(f,[a, b], y)<\infty$, i.e., for a.e. $y \in \mathbb{R}$ because $N(f,[a, b], \cdot)$ is integrable. Thus

$$
\begin{align*}
0 \leq V_{f}(x)-V_{f}\left(x_{0}\right)= & \int_{\mathbb{R}} N(f,[a, x], y) d \lambda(y) \\
& -\int_{\mathbb{R}} N\left(f,\left[a, x_{0}\right], y\right) d \lambda(y) \\
= & \int_{\mathbb{R}} N\left(f,\left(x_{0}, x\right], y\right) d \lambda(y)  \tag{7.48}\\
\leq & \int_{f\left(\left(x_{0}, x\right]\right)} N(f,[a, b], y) d \lambda(y)
\end{align*}
$$

Since $f$ is continuous at $x_{0}, \lambda\left[f\left(\left(x_{0}, x\right]\right)\right] \rightarrow 0$ as $x \rightarrow x_{0}$ and thus

$$
\lim _{x \rightarrow x_{0}^{+}} V_{f}(x)=V_{f}\left(x_{0}\right)
$$

A similar argument shows that

$$
\lim _{x \rightarrow x_{0}^{-}} V_{f}(x)=V_{f}\left(x_{0}\right)
$$

and thus that $V_{f}$ is continuous at $x_{0}$.
Now assume that $f$ also satisfies condition $N$ and let $A$ be a set with $\lambda(A)=0$ so that $\lambda(f(A))=0$. Observe that for any measurable set $E \subset \mathbb{R}$, the set function

$$
\mu(E):=\int_{E} N(f,[a, b], y) d y
$$

is a measure that is absolutely continuous with respect to Lebesgue measure. Consequently, for each $\varepsilon>0$ there exists $\delta>0$ such that $\mu(E)<\varepsilon$ whenever $\lambda(E)<\delta$. Hence if $U \supset f(A)$ is an open set with $\lambda(f(A))<\delta$, we have

$$
\begin{equation*}
\int_{U} N(f,[a, b], y) d y<\varepsilon \tag{7.49}
\end{equation*}
$$

The open set $f^{-1}(U)$ can be expressed as the countable disjoint union of intervals $\bigcup_{i=1}^{\infty} I_{i} \supset A$. Thus, with the notation $I_{i}:=\left[a_{i}, b_{i}\right]$ we have

$$
\begin{array}{rlr}
\lambda\left(V_{f}(A)\right) & \leq \lambda\left(V_{f}\left(\bigcup_{i}^{\infty} I_{i}\right)\right) \\
& \leq \sum_{i=1}^{\infty} \lambda\left(V_{f}\left(I_{i}\right)\right) & \\
& =\sum_{i=1}^{\infty} V_{f}\left(b_{i}\right)-V_{f}\left(a_{i}\right) & \\
& =\sum_{i=1}^{\infty} \int_{\mathbb{R}} N\left(f,\left(a_{i}, b_{i}\right], y\right) d y & \text { by }(7.48) \\
& =\int_{\cup_{i=1}^{\infty} f\left(\left(a_{i}, b_{i}\right]\right)} N(f,[a, b], y) d y & \text { since the } I_{i}^{\prime} s \text { are disjoint } \\
& =\int_{U} N(f,[a, b], y), & \\
& <\varepsilon & \text { by }(7.49)
\end{array}
$$

which implies that $\lambda\left(V_{f}(A)\right)=0$.
7.35. Theorem. If $f$ is a nondecreasing function defined on $[a, b]$, then it is absolutely continuous if and only if the following two conditions are satisfied:
(i) $f$ is continuous,
(ii) $f$ satisfies condition $N$.

Proof. Clearly condition (i) is necessary for absolute continuity and Theorem 7.27 shows that condition (ii) is also necessary.

To prove that the two conditions are sufficient, let $g(x):=x+f(x)$, so that $g$ is a strictly increasing, continuous function. Note that $g$ satisfies condition $N$ since $f$ does and $\lambda(g(I))=\lambda(I)+\lambda(f(I))$ for any interval $I$. Also, if $E$ is a measurable set, then so is $g(E)$ because $E=F \cup N$ where $F$ is an $F_{\sigma}$ set and $N$ has measure zero, by Theorem 4.25. Since $g$ is continuous and since $F$ can be expressed as the countable union of compact sets, it follows that $g(E)=g(F) \cup g(N)$ is the union of a countable number of compacts sets and a set of measure zero and is therefore
a measurable set. Now define a measure $\mu$ by

$$
\mu(E)=\lambda(g(E))
$$

for any measurable set $E$. Observe that $\mu$ is, in fact, a measure since $g$ is injective. Furthermore, $\mu \ll \lambda$ because $g$ satisfies condition $N$. Consequently, the RadonNikodym Theorem applies (Theorem 6.43) and we obtain a function $h \in L^{1}(\lambda)$ such that

$$
\mu(E)=\int_{E} h d \lambda \quad \text { for every measurable set } E .
$$

In particular, taking $E=[a, x]$, we obtain

$$
g(x)-g(a)=\lambda(g(E))=\mu(E)=\int_{E} h d \lambda=\int_{a}^{x} h d \lambda
$$

Thus, as in (7.38), we conclude that $g$ is absolutely continuous and therefore, so is $f$.
7.36. Corollary. A function $f$ defined on $[a, b]$ is absolutely continuous if and only if $f$ satisfies the following three conditions on $[a, b]$ :
(i) $f$ is continuous,
(ii) $f$ is of bounded variation,
(iii) $f$ satisfies condition $N$.

Proof. The necessity of the three conditions is established by Theorems 7.25 and 7.27.

To prove sufficiency suppose $f$ satisfies conditions (i)-(iii). It follows from Corollary 7.34 that $V_{f}(\cdot)$ is continuous, satisfies condition $N$ and therefore is absolutely continuous by Theorem 7.36. Since $f=f_{1}-f_{2}$ where $f_{1}=\frac{1}{2}\left(V_{f}+f\right)$ and $f_{1}=\frac{1}{2}\left(V_{f}-f\right)$, it follows that both $f_{1}$ and $f_{2}$ are absolutely continuous and therefore so is $f$.

### 7.7. Curve Length

Adapting the methods of the previous section, the notion of length is developed and shown to be closely related to 1-dimensional Hausdorff measure.
7.37. Definitions. A curve in $\mathbb{R}^{n}$ is a continuous mapping $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$, and its length is defined as

$$
\begin{equation*}
L_{\gamma}=\sup \sum_{i=1}^{k}\left|\gamma\left(t_{i}\right)-\gamma\left(t_{i-1}\right)\right| \tag{7.50}
\end{equation*}
$$

where the supremum is taken over all finite sequences $a=t_{0}<t_{1}<\cdots<t_{k}=b$. Note that $\gamma(x)$ is a vector in $\mathbb{R}^{n}$ for each $x \in[a, b]$; writing $\gamma(x)$ in terms of its component functions we have

$$
\gamma(x)=\left(\gamma_{1}(x), \gamma_{2}(x), \ldots, \gamma_{n}(x)\right)
$$

Thus, in (7.50) $\gamma\left(t_{i}\right)-\gamma\left(t_{i-1}\right)$ is a vector in $\mathbb{R}^{n}$ and $\left|\gamma\left(t_{i}\right)-\gamma\left(t_{i-1}\right)\right|$ is its length. For $x \in[a, b]$, we will use the notation $L_{\gamma}(x)$ to denote the length of $\gamma$ restricted to the interval $[a, x] . \gamma$ is said to have finite length or to be rectifiable if $L_{\gamma}(b)<\infty$.

We will show that there is a strong parallel between the notions of length and bounded variation. In case $\gamma$ is a curve in $\mathbb{R}$, i.e, in case $\gamma:[a, b] \rightarrow \mathbb{R}$, the two notions coincide. More generally, we have the following.
7.38. Theorem. A continuous curve $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ is rectifiable if and only if each component function, $\gamma_{i}$, is of bounded variation on $[a, b]$.

Proof. Suppose each component function is of bounded variation. Then there are numbers $M_{1}, M_{2}, \cdots, M_{n}$ such that for any finite partition $\mathcal{P}$ of $[a, b]$ into nonoverlapping intervals $I=\left[a_{I}, b_{I}\right]$,

$$
\sum_{I \in \mathcal{P}}\left|\gamma_{j}\left(b_{I}\right)-\gamma_{j}\left(a_{I}\right)\right| \leq M_{j}, j=1,2, \ldots, n
$$

Thus, with $M=M_{1}+M_{2}+\cdots+M_{n}$ we have

$$
\begin{aligned}
\sum_{I \in \mathcal{P}} & \left|\gamma\left(b_{I}\right)-\gamma\left(a_{I}\right)\right| \\
= & \sum_{I \in \mathcal{P}}\left[\left(\gamma_{1}\left(b_{I}\right)-\gamma_{1}\left(a_{I}\right)\right)^{2}+\left(\gamma_{2}\left(b_{I}\right)-\gamma_{2}\left(a_{I}\right)\right)^{2}\right. \\
& \left.+\cdots+\left(\gamma_{n}\left(b_{I}\right)-\gamma_{n}\left(a_{I}\right)\right)^{2}\right]^{1 / 2} \\
\leq & \sum_{I \in \mathcal{P}}\left|\gamma_{1}\left(b_{I}\right)-\gamma_{1}\left(a_{I}\right)\right|+\sum_{I \in \mathcal{P}}\left|\gamma_{2}\left(b_{I}\right)-\gamma_{2}\left(a_{I}\right)\right| \\
\quad & +\cdots+\sum_{I \in \mathcal{P}}\left|\gamma_{n}\left(b_{I}\right)-\gamma_{n}\left(a_{I}\right)\right| \\
\leq & M .
\end{aligned}
$$

Since the partition $\mathcal{P}$ is arbitrary, we conclude that the length of $\gamma$ is less than or equal to $M$.

Now assume that $\gamma$ is rectifiable. Then for any partition $\mathcal{P}$ of $[a, b]$ and any integer $j \in[1, n]$,

$$
\sum_{I \in \mathcal{P}}\left|\gamma\left(b_{I}\right)-\gamma\left(a_{I}\right)\right| \geq \sum_{I \in \mathcal{P}}\left|\gamma_{j}\left(b_{I}\right)-\gamma_{j}\left(a_{I}\right)\right|
$$

thus showing that the total variation of $\gamma_{j}$ is no more than the length of $\gamma$.

In elementary calculus, we know that the formula for the length of a curve $\gamma=\left(\gamma_{1}, \gamma_{2}\right)$ defined on $[a, b]$ is given by

$$
L_{\gamma}=\int_{a}^{b} \sqrt{\left(\gamma_{1}^{\prime}(t)\right)^{2}+\left(\gamma_{2}^{\prime}(t)\right)^{2}} d t
$$

We will proceed to investigate the conditions under which this formula holds using our definition of length. We will consider a curve $\gamma$ in $\mathbb{R}^{n}$; thus we have $\gamma:[a, b] \rightarrow$ $\mathbb{R}^{n}$ with $\gamma(t)=\left(\gamma_{1}(t), \gamma_{2}(t), \ldots, \gamma_{n}(t)\right)$ and we recall the notation $L_{\gamma}(t)$ introduced earlier that denotes the length of $\gamma$ from $a$ to $t$.
7.39. Theorem. If $\gamma$ is rectifiable, then

$$
L_{\gamma}^{\prime}(t)=\sqrt{\left(\gamma_{1}^{\prime}(t)\right)^{2}+\left(\gamma_{2}^{\prime}(t)\right)^{2}+\cdots+\left(\gamma_{n}^{\prime}(t)\right)^{2}}
$$

for almost all $t \in[a, b]$.
Proof. The number $\left|L_{\gamma}(t+h)-L_{\gamma}(t)\right|$ denotes the length along the curve between the points $\gamma(t+h)$ and $\gamma(t)$, which is clearly not less than the straight-line distance. Therefore, it is intuitively clear that

$$
\begin{equation*}
\left|L_{\gamma}(t+h)-L_{\gamma}(t)\right| \geq|\gamma(t+h)-\gamma(t)| \tag{7.51}
\end{equation*}
$$

The rigorous argument to establish this is very similar to the proof of (7.27). Thus, for $h>0$, consider an arbitrary partition $a=t_{0}<t_{1}<\cdots<t_{k}=x$. Then from the definition of $L_{\gamma}$,

$$
L_{\gamma}(x+h) \geq|\gamma(x+h)-\gamma(x)|+\sum_{i=1}^{k}\left|\gamma\left(t_{i}\right)-\gamma\left(t_{i-1}\right)\right|
$$

Since

$$
L_{\gamma}(x)=\sup \left\{\sum_{i=1}^{k}\left|\gamma\left(t_{i}\right)-\gamma\left(t_{i-1}\right)\right|\right\}
$$

over all partitions $a=t_{0}<t_{1}<\cdots<t_{k}=x$, it follows

$$
L_{\gamma}(x+h) \geq|\gamma(x+h)-\gamma(x)|+L_{\gamma}(x)
$$

A similar inequality holds for $h<0$ and therefore we obtain

$$
\begin{equation*}
\left|L_{\gamma}(t+h)-L_{\gamma}(t)\right| \geq|\gamma(t+h)-\gamma(t)| \tag{7.52}
\end{equation*}
$$

consequently,

$$
\begin{align*}
\left|L_{\gamma}(t+h)-L_{\gamma}(t)\right| \geq & |\gamma(t+h)-\gamma(t)| \\
= & {\left[\left(\gamma_{1}(t+h)-\gamma_{1}(t)\right)^{2}+\left(\gamma_{2}(t+h)-\gamma_{2}(t)\right)^{2}\right.}  \tag{7.53}\\
& \left.\quad+\cdots+\left(\gamma_{n}(t+h)-\gamma_{n}(t)\right)^{2}\right]^{1 / 2}
\end{align*}
$$

whenever $t+h, t \in[a, b]$. Consequently,

$$
\begin{aligned}
\frac{\left|L_{\gamma}(t+h)-L_{\gamma}(t)\right|}{h} \geq & {\left[\left(\frac{\gamma_{1}(t+h)-\gamma_{1}(t)}{h}\right)^{2}\right.} \\
& \left.+\cdots+\left(\frac{\gamma_{n}(t+h)-\gamma_{n}(t)}{h}\right)^{2}\right]^{1 / 2}
\end{aligned}
$$

Taking the limit as $h \rightarrow 0$, we obtain

$$
\begin{equation*}
L_{\gamma}^{\prime}(t) \geq \sqrt{\left(\gamma_{1}^{\prime}(t)\right)^{2}+\left(\gamma_{2}^{\prime}(t)\right)^{2}+\cdots+\left(\gamma_{n}^{\prime}(t)\right)^{2}} \tag{7.54}
\end{equation*}
$$

whenever all derivatives exist, which is almost everywhere in view of Theorem 7.38. The remainder of the proof is completely analogous to the proof of (7.29) in Theorem 7.22 and is left as an exercise.

The proof of the following result is completely analogous to the proof of Theorem 7.32 and is also left as an exercise.
7.40. TheOrem. If each component function $\gamma_{i}$ of the curve $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ is absolutely continuous, the function $L_{\gamma}(\cdot)$ is absolutely continuous and

$$
L_{\gamma}(x)=\int_{a}^{x} L_{\gamma}^{\prime} d \lambda=\int_{a}^{x} \sqrt{\left(\gamma_{1}^{\prime}\right)^{2}+\left(\gamma_{2}^{\prime}\right)^{2}+\cdots+\left(\gamma_{n}^{\prime}\right)^{2}} d \lambda
$$

for each $x \in[a, b]$.
7.41. Example. Intuitively, one might expect that the trace of a continuous curve would resemble a piece of string, perhaps badly crumpled, but still like a piece of string. However, Peano discovered that the situation could be far worse. He was the first to demonstrate the existence of a continuous mapping $\gamma:[0,1] \rightarrow \mathbb{R}^{2}$ such that $\gamma[0,1]$ occupies the unit square, $Q$. In other words, he showed the existence of an "area-filling curve." In the figure below, we show the first three stages of the construction of such a curve. This construction is due to Hilbert.


Each stage represents the graph of a continuous (piecewise linear) mapping $\gamma_{k}:[0,1] \rightarrow \mathbb{R}^{2}$. From the way the construction is made, we find that

$$
\sup _{t \in[0,1]}\left|\gamma_{k}(t)-\gamma_{l}(t)\right| \leq \frac{\sqrt{2}}{2^{k}}
$$

where $k \leq l$. Hence, since the space of continuous functions with the topology of uniform convergence is complete, there exists a continuous mapping $\gamma$ that is the uniform limit of $\left\{\gamma_{k}\right\}$. To see that $\gamma[0,1]$ is the unit square $Q$, first observe that each point $x_{0}$ in the unit square belongs to some square of each of the partitions of $Q$. Denoting by $\mathcal{P}_{k}$ the $k^{\text {th }}$ partition of $Q$ into $4^{k}$ subsquares of side-length $(1 / 2)^{k}$, we see that each point $x_{0} \in Q$ belongs to some square, $Q_{k}$, of $\mathcal{P}_{k}$ for $k=1,2, \ldots$. For each $k$ the curve $\gamma_{k}$ passes through $Q_{k}$, and so it is clear that there exist points $t_{k} \in[0,1]$ such that

$$
\begin{equation*}
x_{0}=\lim _{k \rightarrow \infty} \gamma_{k}\left(t_{k}\right) \tag{7.55}
\end{equation*}
$$

For $\varepsilon>0$, choose $K_{1}$ such that $\left|x_{0}-\gamma_{k}\left(t_{k}\right)\right|<\varepsilon / 2$ for $k \geq K_{1}$. Since $\gamma_{k} \rightarrow \gamma$ uniformly, we see by Theorem 3.56, that there exists $\delta=\delta(\varepsilon)>0$ such that $\left|\gamma_{k}(t)-\gamma_{k}(s)\right|<\varepsilon$ for all $k$ whenever $|s-t|<\delta$. Choose $K_{2}$ such that $\left|t-t_{k}\right|<\delta$ for $k \geq K_{2}$. Since $[0,1]$ is compact, there is a point $t_{0} \in[0,1]$ such that (for a subsequence) $t_{k} \rightarrow t_{0}$ as $k \rightarrow \infty$. Then $x_{0}=\gamma\left(t_{0}\right)$ because

$$
\left|x_{0}-\gamma\left(t_{0}\right)\right| \leq\left|x_{0}-\gamma_{k}\left(t_{k}\right)\right|+\left|\gamma_{k}\left(t_{k}\right)-\gamma_{k}\left(t_{0}\right)\right|<\varepsilon
$$

for $k \geq \max K_{1}, K_{2}$. One must be careful to distinguish a curve in $\mathbb{R}^{n}$ from the point set described by its trace. For example, compare the curve in $\mathbb{R}^{2}$ given by

$$
\gamma(x)=(\cos x, \sin x), x \in[0,2 \pi]
$$

to the curve

$$
\eta(x)=(\cos 2 x, \sin 2 x), x \in[0,2 \pi]
$$

Their traces occupy the same point set, namely, the unit circle. However, the length of $\gamma$ is $2 \pi$ whereas the length of $\eta$ is $4 \pi$. This simple example serves as a model for the relationship between the length of a curve and the 1-dimensional Hausdorff measure of its trace. Roughly speaking, we will show that they are the same if one takes into account the number of times each point in the trace is covered. In particular, if $\gamma$ is injective, then they are the same.
7.42. Theorem. Let $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ be a continuous curve. Then

$$
\begin{equation*}
L_{\gamma}(b)=\int N(\gamma,[a, b], y) d H^{1}(y) \tag{7.56}
\end{equation*}
$$

where $N(\gamma,[a, b], y)$ denotes the (possibly infinite) number of points in the set $\gamma^{-1}\{y\} \cap[a, b]$. Equality in (7.56) is understood in the sense that either both sides are finite and equal or both sides are infinite.

Proof. Let $M$ denote the length of $\gamma$ on $[a, b]$. We will first prove

$$
\begin{equation*}
M \geq \int N(\gamma,[a, b], y) d H^{1}(y) \tag{7.57}
\end{equation*}
$$

and therefore, we may as well assume $M<\infty$. The function, $L_{\gamma}(\cdot)$, is nondecreasing, and its range is the interval $[0, M]$. As in (7.53), we have

$$
\begin{equation*}
\left|L_{\gamma}(x)-L_{\gamma}(y)\right| \geq|\gamma(x)-\gamma(y)| \tag{7.58}
\end{equation*}
$$

for all $x, y \in[a, b]$. Since $L_{\gamma}$ is nondecreasing, it follows that $L_{\gamma}^{-1}\{s\}$ is an interval (possibly degenerate) for all $s \in[0, M]$. In fact, there are only countably many $s$ for which $L_{\gamma}^{-1}\{s\}$ is a nondegenerate interval. Now define $g:[0, M] \rightarrow \mathbb{R}^{n}$ as follows:

$$
\begin{equation*}
g(s)=\gamma(x) \tag{7.59}
\end{equation*}
$$

where $x$ is any point in $L_{\gamma}^{-1}(s)$. Observe that if $L_{\gamma}^{-1}\{s\}$ is a nondegenerate interval and if $x_{1}, x \in L_{\gamma}^{-1}\{s\}$, then $\gamma\left(x_{1}\right)=\gamma(x)$, thus ensuring that $g(s)$ is well-defined. Also, for $x \in L_{\gamma}^{-1}\{s\}, y \in L_{\gamma}^{-1}\{t\}$, notice from (7.58) that

$$
\begin{equation*}
|g(s)-g(t)|=|\gamma(x)-\gamma(y)| \leq\left|L_{\gamma}(x)-L_{\gamma}(y)\right|=|s-t| \tag{7.60}
\end{equation*}
$$

so that $g$ is a Lipschitz function with Lipschitz constant 1. From the definition of $g$, we clearly have $\gamma=g \circ L_{\gamma}$. Let $S$ be the set of points in $[0, M]$ such that $L_{\gamma}^{-1}\{s\}$ is a nondegenerate interval. Then it follows that

$$
N(\gamma,[a, b], y)=N(g,[0, M], y)
$$

for all $y \notin g(S)$. Since $S$ is countable and therefore also $g(S)$, we have

$$
\begin{equation*}
\int N(\gamma,[a, b], y) d H^{1}(y)=\int N(g,[0, M], y) d H^{1}(y) \tag{7.61}
\end{equation*}
$$

We will appeal to the proof of Theorem 7.33 to show that

$$
M \geq \int N(g,[0, M], y) d H^{1}(y)
$$

Let $\mathcal{P}_{i}$ be a sequence of partitions of $[0, M]$ each having the property that its intervals have length less than $1 / i$. Then, using the argument that established (7.44), we have

$$
N(g,[0, M], y)=\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \chi_{g(I)}(y)\right\}
$$

Since

$$
H^{1}(g(I))=\int \chi_{g(I)}(y) d H^{1}(y)
$$

for each interval $I$, we adapt the proof of (7.46) to obtain

$$
\begin{equation*}
\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} H^{1}[g(I)]\right\}=\int N(g,[0, M], y) d H^{1}(y) \tag{7.62}
\end{equation*}
$$

Now use (7.60) and Exercise 2, Section 7.7, to conclude that $H^{1}[g(I)] \leq \lambda(I)$ so that (7.62) yields

$$
\begin{equation*}
\liminf _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} \lambda(I)\right\} \geq \int_{0}^{\infty} N(g,[0, M], y) d H^{1}(y) \tag{7.63}
\end{equation*}
$$

It is necessary to use the liminf here because we don't know that the limit exists.
However, since

$$
\sum_{I \in \mathcal{P}_{i}} \lambda(I)=\lambda([0, M])=M
$$

we see that the limit does exist and that the left-side of (7.63) equals $M$. Thus, we obtain

$$
M \geq \int_{0}^{\infty} N(g,[0, M], y) d H^{1}(y)
$$

which, along with (7.61), establishes (7.57).
We now will prove

$$
\begin{equation*}
M \leq \int_{0}^{\infty} N(\gamma,[a, b], y) d H^{1}(y) \tag{7.64}
\end{equation*}
$$

to conclude the proof of the theorem. Again, we adapt the reasoning leading to (7.62) to obtain

$$
\begin{equation*}
\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}} H^{1}[\gamma(I)]\right\}=\int N(\gamma,[a, b], y) d H^{1}(y) \tag{7.65}
\end{equation*}
$$

In this context, $\mathcal{P}_{i}$ is a sequence of partitions of $[a, b]$, each of whose intervals has maximum length $1 / i$. Each term on the left-side involves $H^{1}[\gamma(I)]$. Now $\gamma(I)$ is the trace of a curve defined on the interval $I=\left[a_{I}, b_{I}\right]$. By Exercise 2, Section 7.7, we know that $H^{1}[\gamma(I)]$ is greater than or equal to the $H^{1}$ measure of the orthogonal projection of $\gamma(I)$ onto any straight line, $l$. That is, if $p: \mathbb{R}^{n} \rightarrow l$ is an orthogonal projection, then

$$
H^{1}[p(\gamma(I))] \leq H^{1}[\gamma(I)]
$$

In particular, consider the straight line $l$ that passes through the points $\gamma\left(a_{I}\right)$ and $\gamma\left(b_{I}\right)$. Since $I$ is connected and $\gamma$ is continuous, $\gamma(I)$ is connected and therefore, its projection, $p[\gamma(I)]$, onto $l$ is also connected. Thus, $p[\gamma(I)]$ must contain the interval with endpoints $\gamma\left(a_{I}\right)$ and $\gamma\left(b_{I}\right)$. Hence $\left|\gamma\left(b_{I}\right)-\gamma\left(a_{I}\right)\right| \leq H^{1}[\gamma(I)]$. Thus, from (7.65), we have

$$
\lim _{i \rightarrow \infty}\left\{\sum_{I \in \mathcal{P}_{i}}\left|\gamma\left(b_{I}\right)-\gamma\left(a_{I}\right)\right|\right\} \leq \int N(\gamma,[a, b], y) d H^{1}(y)
$$

By Exercise 7.3 the expression on the left is the length of $\gamma$ on $[a, b]$, which is $M$, thus proving (7.64).
7.43. Remark. The function $g$ defined in (7.59) is called a para-metrization of $\gamma$ with respect to arc-length. The purpose of $g$ is to give an alternate and equivalent description of the curve $\gamma$. It is equivalent in the sense that the trace and length of $g$ are the same as those of $\gamma$. In case $\gamma$ is not constant on any interval, then $L_{\gamma}$ is a homeomorphism and thus $g$ and $\gamma$ are related by a homeomorphic change of variables.

## Exercises for Section 7.7

1. In the proof of Theorem 7.48 we established (7.69) by means of Lemma 7.47. Another way to obtain (7.69) is the following. Prove that if $f$ is Lipschitz on a set $E$ with Lipschitz constant $C$, then $\lambda[f(A)] \leq C \lambda(A)$ whenever $A \subset E$ is a Lebesgue measurable. This is the same as Exercise 4, Section 7.5, except that $f$ is defined only on $E$, not necessarily on $\mathbb{R}$. First prove that Lebesgue measure can be defined as follows:

$$
\lambda(A)=\inf \left\{\sum_{i=1}^{\infty} \operatorname{diam} E_{i}: A \subset \bigcup_{i=1}^{\infty} E_{i}\right\}
$$

where the $E_{i}$ are arbitrary sets.
2. Suppose $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is a Lipschitz mapping with Lipschitz constant $C$. Prove that

$$
H^{k}[f(E)] \leq C^{k} H^{k}(E)
$$

for $E \subset \mathbb{R}^{n}$.
3. Prove that if $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ is a continuous curve and $\left\{\mathcal{P}_{i}\right\}$ is a sequence of partitions of $[a, b]$ such that

$$
\lim _{i \rightarrow \infty} \max _{I \in \mathcal{P}_{i}}\left|b_{I}-a_{I}\right|=0
$$

then

$$
\lim _{i \rightarrow \infty} \sum_{I \in \mathcal{P}_{i}}\left|\gamma\left(b_{I}\right)-\gamma\left(a_{I}\right)\right|=L_{\gamma}(b)
$$

4. Prove that the example of an area filling curve in Section 7.7 actually has the unit square as its trace.
5. It follows from Theorem 7.42 that the area filling curve is not rectifiable. Prove this directly from the construction of the curve.
6. Give an example of a continuous curve that fills the unit cube in $\mathbb{R}^{3}$.
7. Give a proof of Lemma 7.40.
8. Let $\gamma:[0,1] \rightarrow \mathbb{R}^{2}$ be defined by $\gamma(x)=(x, f(x))$ where $f$ is the Cantor-Lebesgue function described in Example 5.7. Thus, $\gamma$ describes the graph of the Cantor function. Find the length of $\gamma$.

### 7.8. The Critical Set of a Function

During the course of our development of the Fundamental Theorem of Calculus in Section 7.5, we found that absolutely continuous functions are continuous functions of bounded variation that satisfy condition $N$. We will show here that these properties characterize AC functions. This will be done by carefully analyzing the behavior of a function on the set where its derivative is 0 .

Recall that a function $f$ defined on $[a, b]$ is said to satisfy condition $N$ provided $\lambda[f(E)]=0$ whenever $\lambda(E)=0$ for $E \subset[a, b]$.

An example of a function that does not satisfy condition $N$ is the CantorLebesgue function. Indeed, it maps the Cantor set (of Lebesgue measure zero) onto the unit interval $[0,1]$. On the other hand, a Lipschitz function is an example of a function that does satisfy condition $N$ (see Exercise 3, Section 7.5). Recall Definition 3.10, which states that $f$ satisfies a Lipschitz condition on [a,b] if there exists a constant $C=C_{f}$ such that

$$
|f(x)-f(y)| \leq C|x-y| \text { whenever } x, y \in[a, b]
$$

One of the important aspects of a function is its behavior on the critical set, the set where its derivative is zero. One would expect that the critical set of a function $f$ would be mapped onto a set of measure zero since $f$ is neither increasing nor decreasing at points where $f^{\prime}=0$. For convenience, we state this result which was proved earlier, Theorem 7.29.
7.44. Theorem. Suppose $f$ is defined on $[a, b]$. Let

$$
E=(a, b) \cap\left\{x: f^{\prime}(x)=0\right\} .
$$

Then $\lambda[f(E)]=0$.
Now we will investigate the behavior of a function on the complement of its critical set and show that good things happen there. We will prove that the set on which a continuous function has a non-zero derivative can be decomposed into a countable collection of disjoint sets on each of which the function is bi-Lipschitzian. That is, on each of these sets the function is Lipschitz and injective; furthermore, its inverse is Lipschitz on the image of each such set.
7.45. Theorem. Suppose $f$ is defined on $[a, b]$ and let $A$ be defined by

$$
A:=[a, b] \cap\left\{x: f^{\prime}(x) \text { exists and } f^{\prime}(x) \neq 0\right\}
$$

Then for each $\theta>1$, there is a countable collection $\left\{E_{k}\right\}$ of disjoint Borel sets such that
(i) $A=\bigcup_{k=1}^{\infty} E_{k}$
(ii) For each positive integer $k$ there is a positive rational number $r_{k}$ such that

$$
\begin{align*}
\frac{r_{k}}{\theta} & \leq\left|f^{\prime}(x)\right| \leq \theta r_{k} \quad \text { for } \quad x \in E_{k} \\
|f(y)-f(x)| & \leq \theta r_{k}|x-y| \quad \text { for } \quad x, y \in E_{k}  \tag{7.66}\\
|f(y)-f(x)| & \geq \frac{r_{k}}{\theta}(y-x) \quad \text { for } \quad x, y \in E_{k} \tag{7.67}
\end{align*}
$$

Proof. Since $\theta>1$ there exists $\varepsilon>0$ so that

$$
\frac{1}{\theta}+\varepsilon<1<\theta-\varepsilon
$$

For each positive integer $k$ and each positive rational number $r$ let $A(k, r)$ be the set of all points $x \in A$ such that

$$
\left(\frac{1}{\theta}+\varepsilon\right) r \leq\left|f^{\prime}(x)\right| \leq(\theta-\varepsilon) r
$$

With the help of the triangle inequality, observe that if $x, y \in[a, b]$ with $x \in A(k, r)$ and $|y-x|<1 / k$, then

$$
|f(y)-f(x)| \leq \varepsilon r|(y-x)|+\left|f^{\prime}(x)(y-x)\right| \leq \theta r|y-x|
$$

and

$$
|f(y)-f(x)| \geq-\varepsilon r|y-x|+\left|f^{\prime}(x)(y-x)\right| \geq \frac{r}{\theta}|y-x|
$$

Clearly,

$$
A=\bigcup A(k, r)
$$

where the union is taken as $k$ and $r$ range through the positive integers and positive rationals, respectively. To ensure that (7.66) and (7.67) hold for all $a, b \in E_{k}$ (defined below), we express $A(k, r)$ as the countable union of sets each having diameter $1 / k$ by writing

$$
A(k, r)=\bigcup_{s \in \mathbb{Q}}[A(k, r) \cap I(s, 1 / k)]
$$

where $I(s, 1 / k)$ denotes the open interval of length $1 / k$ centered at the rational number $s$. The sets $[A(k, r) \cap I(s, 1 / k)]$ constitute a countable collection as $k, r$,
and $s$ range through their respective sets. Relabel the sets $[A(k, r) \cap I(s, 1 / k)]$ as $E_{k}, k=1,2, \ldots$. We may assume the $E_{k}$ are disjoint by appealing to Lemma 4.7, thus obtaining the desired result.

The next result differs from the preceding one only in that the hypothesis now allows the set $A$ to include critical points of $f$. There is no essential difference in the proof.
7.46. Theorem. Suppose $f$ is defined on $[a, b]$ and let $A$ be defined by

$$
A=[a, b] \cap\left\{x: f^{\prime}(x) \text { exists }\right\}
$$

Then for each $\theta>1$, there is a countable collection $\left\{E_{k}\right\}$ of disjoint sets such that
(i) $A=\cup_{k=1}^{\infty} E_{k}$,
(ii) For each positive integer $k$ there is a positive rational number $r$ such that

$$
\left|f^{\prime}(x)\right| \leq \theta r \quad \text { for } \quad x \in E_{k}
$$

and

$$
|f(y)-f(x)| \leq \theta r|y-x| \quad \text { for } \quad x, y \in E_{k}
$$

Before giving the proof of the next main result, we take a slight diversion that is concerned with the extension of Lipschitz functions. We will give a proof of a special case of Kirzbraun's Theorem whose general formulation we do not require and is more difficult to prove.
7.47. Theorem. Let $A \subset \mathbb{R}$ be an arbitrary set and suppose $f: A \rightarrow \mathbb{R}$ is a Lipschitz function with Lipschitz constant $C$. Then there exists a Lipschitz function $\bar{f}: \mathbb{R} \rightarrow \mathbb{R}$ with the same Lipschitz constant $C$ such that $\bar{f}=f$ on $A$.

Proof. Define

$$
\bar{f}(x)=\inf \{f(a)+C|x-a|: a \in A\}
$$

Clearly, $\bar{f}=f$ on $A$ because if $b \in A$, then

$$
f(b)-f(a) \leq C|b-a|
$$

for any $a \in A$. This shows that $f(b) \leq \bar{f}(b)$. On the other hand, $\bar{f}(b) \leq f(b)$ follows immediately from the definition of $\bar{f}$. Finally, to show that $\bar{f}$ has Lipschitz
constant $C$, let $x, y \in \mathbb{R}$. Then

$$
\begin{aligned}
\bar{f}(x) & \leq \inf \{f(a)+C(|y-a|+|x-y|): a \in A\} \\
& =\bar{f}(y)+C|x-y|,
\end{aligned}
$$

which proves $\bar{f}(x)-\bar{f}(y) \leq C|x-y|$. The proof with $x$ and $y$ interchanged is similar.
7.48. Theorem. Suppose $f$ is a continuous function on $[a, b]$ and let

$$
A=[a, b] \cap\left\{x: f^{\prime}(x) \text { exists and } f^{\prime}(x) \neq 0\right\} .
$$

Then, for every Lebesgue measurable set $E \subset A$,

$$
\begin{equation*}
\int_{E}\left|f^{\prime}\right| d \lambda=\int_{-\infty}^{\infty} N(f, E, y) d \lambda(y) \tag{7.68}
\end{equation*}
$$

Equality is understood in the sense that either both sides are finite and equal or both sides are infinite.

Proof. We apply Theorem 7.45 with $A$ replaced by $E$. Thus, for each $k \in \mathbb{N}$ there is a positive rational number $r$ such that

$$
\frac{r}{\bar{\theta}} \lambda\left(E_{k}\right) \leq \int_{E_{k}}\left|f^{\prime}\right| d \lambda \leq \theta r \lambda\left(E_{k}\right)
$$

and $f$ restricted to $E_{k}$ satisfies a Lipschitz condition with constant $\theta r$. Therefore by Lemma $7.47, f$ has a Lipschitz extension to $\mathbb{R}$ with the same Lipschitz constant. From Exercise 4, Section 7.5, we obtain

$$
\begin{equation*}
\lambda\left[f\left(E_{k}\right)\right] \leq \theta r \lambda\left(E_{k}\right) \tag{7.69}
\end{equation*}
$$

Theorem 7.45 states that $f$ restricted to $E_{k}$ is univalent and that its inverse function is Lipschitz with constant $\theta / r$. Thus, with the same reasoning as before,

$$
\lambda\left(E_{k}\right) \leq \frac{\theta}{r} \lambda\left[f\left(E_{k}\right)\right] .
$$

Hence, we obtain

$$
\begin{equation*}
\frac{1}{\theta^{2}} \lambda\left(f\left[E_{k}\right]\right) \leq \int_{E_{k}}\left|f^{\prime}\right| d \lambda \leq \theta^{2} \lambda\left[f\left(E_{k}\right)\right] . \tag{7.70}
\end{equation*}
$$

Each $E_{k}$ can be expressed as the countable union of compact sets and a set of Lebesgue measure zero. Since $f$ restricted to $E_{k}$ satisfies a Lipschitz condition, $f$ maps the set of measure zero into a set of measure zero and each compact set is
mapped into a compact set. Consequently, $f\left(E_{k}\right)$ is the countable union of compact sets and a set of measure zero and therefore, Lebesgue measurable. Let

$$
g(y)=\sum_{k=1}^{\infty} \chi_{f\left(E_{k}\right)}(y)
$$

so that $g(y)$ is the number of sets $\left\{f\left(E_{k}\right)\right\}$ that contain $y$. Observe that $g$ is Lebesgue measurable. Since the sets $\left\{E_{k}\right\}$ are disjoint, their union is $E$, and $f$ restricted to each $E_{k}$ is univalent, we have

$$
g(y)=N(f, E, y)
$$

Finally,

$$
\frac{1}{\theta^{2}} \sum_{k=1}^{\infty} \lambda\left(f\left[E_{k}\right]\right) \leq \int_{E}\left|f^{\prime}\right| d \lambda \leq \theta^{2} \sum_{k=1}^{\infty} \lambda\left[f\left(E_{k}\right)\right]
$$

and, with the aid of Corollary 6.15

$$
\begin{aligned}
\int N(f, E, y) d \lambda(y) & =\int g(y) d \lambda(y) \\
& =\int \sum_{k=1}^{\infty} \chi_{f\left(E_{k}\right)}(y) d \lambda(y) \\
& =\sum_{k=1}^{\infty} \int \chi_{f\left(E_{k}\right)}(y) d \lambda(y) \\
& =\sum_{k=1}^{\infty} \lambda\left[f\left(E_{k}\right)\right]
\end{aligned}
$$

The result now follows from (7.70) since $\theta>1$ is arbitrary.
7.49. Corollary. If $f$ satisfies (7.68), then $f$ satisfies condition $N$ on the set A.

We conclude this section with a another result concerning absolute continuity. Theorem 7.46 states that a function possesses some regularity properties on the set where it is differentiable. Therefore, it seems reasonable to expect that if $f$ is differentiable everywhere in its domain of definition, then it will have to be a "nice" function. This is the thrust of the next result.
7.50. Theorem. Suppose $f:(a, b) \rightarrow \mathbb{R}$ has the property that $f^{\prime}$ exists everywhere and $f^{\prime}$ is integrable. Then $f$ is absolutely continuous.

Proof. Referring to Theorem 7.45, we find that $(a, b)$ can be written as the union of a countable collection $\left\{E_{k}, k=1,2, \ldots\right\}$ of disjoint Borel sets such that the restriction of $f$ to each $E_{k}$ is Lipschitzian. Hence, it follows that $f$ satisfies
condition $N$ on $(a, b)$. Since $f$ is continuous on $(a, b)$, it remains to show that $f$ is of bounded variation.

For this, let $E_{0}=(a, b) \cap\left\{x: f^{\prime}(x)=0\right\}$. According to Theorem 7.44, $\lambda\left[f\left(E_{0}\right)\right]=0$. Therefore, Theorem 7.48 implies

$$
\int_{E_{k}}\left|f^{\prime}\right| d \lambda=\int_{-\infty}^{\infty} N\left(f, E_{k}, y\right) d \lambda(y)
$$

for $k=1,2, \ldots$. Hence,

$$
\int_{a}^{b}\left|f^{\prime}\right| d \lambda=\int_{-\infty}^{\infty} N(f,(a, b), y) d \lambda(y)
$$

and since $f^{\prime}$ is integrable by assumption, it follows that $f$ is of bounded variation by Theorem 7.33.

## Exercises for Section 7.8

1. Show that the conclusion of Theorem 7.50 still holds if the following assumptions are satisfied: $f^{\prime}$ to exists everywhere on $(a, b)$ except for a countable set, $f^{\prime}$ is integrable, and $f$ is continuous.
2. Prove that the sets $A(k, r)$ defined in the proof of Theorem 7.45 are Borel sets. Hints:
(a) For each positive rational number $r$, let $A_{1}(r)$ denote all points $x \in A$ such that

$$
\left(\frac{1}{\theta}+\varepsilon\right) r \leq\left|f^{\prime}(x)\right| \leq(\theta-\varepsilon) r
$$

Show that $A_{1}(r)$ is a Borel set.
(b) Let $f$ be a continuous function on $[a, b]$. Let

$$
F_{1}(x, y):=\frac{f(y)-f(x)}{x-y} \quad \text { for all } x, y \in[a, b] \text { with } a \neq b
$$

Prove that $F_{1}$ is a Borel function on $[a, b] \times[a, b] \backslash\{(x, y): x=y\}$.
(c) Let $F_{2}(x, y):=f^{\prime}(x)$ for $x \in[a, b]$. Show that $F_{2}$ is a Borel function on $A \times \mathbb{R}$.
(d) For each positive integer $k$, let $A(k)=\{(x, y):|y-x|<1 / k\}$. Note that $A(k)$ is open.
(e) $A(k, r)$ is thus a Borel set since

$$
A(k, r)=\left\{x:\left|F_{1}(x, y)-F_{2}(x, y)\right| \leq \varepsilon r\right\} \cap A_{1}(r) \cap\{x:(x, y) \in A(k)\}
$$

### 7.9. Approximate Continuity

In Section 7.2 the notion of Lebesgue point allowed us to define an integrable function, $f$, at almost all points in a way that does not depend on the choice of any function in the equivalence class determined by $f$. In the development below the concept of approximate continuity will permit us to carry through a similar program for functions that are merely measurable.

A key ingredient in the development of Section 7.2 occurred in the proof of Theorem 7.11, where a continuous function was used to approximate an integrable function in the $L^{1}$-norm. A slightly disquieting feature of this development is that it does not allow the approximation of measurable functions, only integrable ones. In this section this objection is addressed by introducing the concept of approximate continuity.

Throughout this section, we use the following notation. Recall that some of it was introduced in (7.17) and (7.18).

$$
\begin{gathered}
A_{t}=\{x: f(x)>t\} \\
B_{t}=\{x: f(x)<t\} \\
\bar{D}(E, x)=\limsup _{r \rightarrow 0} \frac{\lambda(E \cap B(x, r))}{\lambda(B(x, r))}
\end{gathered}
$$

and

$$
\underline{D}(E, x)=\liminf _{r \rightarrow 0} \frac{\lambda(E \cap B(x, r))}{\lambda(B(x, r))} .
$$

In case the upper and lower limits are equal, we denote their common value by $D(E, x)$. Note that the sets $A_{t}$ and $B_{t}$ are defined up to sets of Lebesgue measure zero.
7.51. Definition. Before giving the next definition, let us first review the definition of limit superior of a function that we discussed earlier on page p. 70. Recall that

$$
\limsup _{x \rightarrow x_{0}} f(x):=\lim _{r \rightarrow 0} M\left(x_{0}, r\right)
$$

where $M\left(x_{0}, r\right)=\sup \left\{f(x): 0<\left|x-x_{0}\right|<r\right\}$. Since $M\left(x_{0}, r\right)$ is a nondecreasing function of $r$, the limit of the right-side exists. If we denote $L\left(x_{0}\right):=$ $\limsup _{x \rightarrow x_{0}} f(x)$, let

$$
\begin{equation*}
T:=\left\{t: B\left(x_{0}, r\right) \cap A_{t}=\emptyset \quad \text { for all small } r\right\} \tag{7.71}
\end{equation*}
$$

If $t \in T$, then there exists $r_{0}>0$ such that for all $0<r<r_{0}, M\left(x_{0}, r\right) \leq$ $t$. Since $M\left(x_{0}, r\right) \downarrow L\left(x_{0}\right)$ it follows that $L\left(x_{0}\right) \leq t$, and therefore, $L\left(x_{0}\right)$ is a lower bound for $T$. On the other hand, if $L\left(x_{0}\right)<t<t^{\prime}$, then the definition of $M\left(x_{0}, r\right)$ implies that $M\left(x_{0}, r\right)<t^{\prime}$ for all small $r>0, \Longrightarrow B\left(x_{0}, r\right) \cap A_{t^{\prime}}=$
$\emptyset$ for all small $r$. As this is true for each $t^{\prime}>L\left(x_{0}\right)$, we conclude that $t$ is not a lower bound for $T$ and therefore that $L\left(x_{0}\right)$ is the greatest lower bound; that is,

$$
\begin{equation*}
\limsup _{x \rightarrow x_{0}} f(x)=L\left(x_{0}\right)=\inf T . \tag{7.72}
\end{equation*}
$$

With the above serving as motivation, we proceed with the measure-theoretic counterpart of (7.72): If $f$ is a Lebesgue measurable function defined on $\mathbb{R}^{n}$, the upper (lower) approximate limit of $f$ at a point $x_{0}$ is defined by

$$
\begin{aligned}
& \operatorname{ap} \limsup _{x \rightarrow x_{0}} f(x)=\inf \left\{t: D\left(A_{t}, x_{0}\right)=0\right\} \\
& \text { ap } \liminf _{x \rightarrow x_{0}} f(x)=\sup \left\{t: D\left(B_{t}, x_{0}\right)=0\right\}
\end{aligned}
$$

We speak of the approximate limit of $f$ at $x_{0}$ when

$$
\operatorname{ap} \limsup _{x \rightarrow x_{0}} f(x)=\text { ap } \liminf _{x \rightarrow x_{0}} f(x)
$$

and $f$ is said to be approximately continuous at $x_{0}$ if

$$
\text { ap } \lim _{x \rightarrow x_{0}} f(x)=f\left(x_{0}\right)
$$

Note that if $g=f$ a.e., then the sets $A_{t}$ and $B_{t}$ corresponding to $g$ differ from those for $f$ by at most a set of Lebesgue measure zero, and thus the upper and lower approximate limits of $g$ coincide with those of $f$ everywhere.

In topology, a point $x$ is interior to a set $E$ if there is a ball $B(x, r) \subset E$. In other words, $x$ is interior to $E$ if it is completely surrounded by other points in $E$. In measure theory, it would be natural to say that $x$ is interior to $E$ (in the measure theoretic sense) if $D(E, x)=1$. See Exercise 1, Section 7.9.

The following is a direct consequence of Theorem 7.11, which implies that almost every point of a measurable set is interior to it (in the measure-theoretic sense).
7.52. Theorem. If $E \subset \mathbb{R}^{n}$ is a Lebesgue measurable set, then

$$
\begin{array}{ll}
D(E, x)=1 & \text { for } \lambda \text {-almost all } x \in E, \\
D(E, x)=0 & \text { for } \lambda \text {-almost all } x \in \widetilde{E} .
\end{array}
$$

Recall that a function $f$ is continuous at $x$ if for every open interval $I$ containing $f(x), x$ is interior to $f^{-1}(I)$. This remains true in the measure-theoretic context.
7.53. Theorem. Suppose $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is Lebesgue measurable function. Then $f$ is approximately continuous at $x$ if and only if for every open interval I containing $f(x), D\left[f^{-1}(I), x\right]=1$.

Proof. Assume $f$ is approximately continuous at $x$ and let $I$ be an arbitrary open interval containing $f(x)$. We will show that $D\left[f^{-1}(I), x\right]=1$. Let $J=\left(t_{1}, t_{2}\right)$ be an interval containing $f(x)$ whose closure is contained in $I$. From the definition of approximate continuity, we have

$$
D\left(A_{t_{2}}, x\right)=D\left(B_{t_{1}}, x\right)=0
$$

and therefore

$$
D\left(A_{t_{2}} \cup B_{t_{1}}, x\right)=0
$$

Since

$$
\mathbb{R}^{n}-f^{-1}(I) \subset A_{t_{2}} \cup B_{t_{1}}
$$

it follows that $D\left(\mathbb{R}^{n}-f^{-1}(I)\right)=0$ and therefore that $D\left[f^{-1}(I), x\right]=1$ as desired.
For the proof in the opposite direction, assume $D\left(f^{-1}(I), x\right)=1$ whenever $I$ is an open interval containing $f(x)$. Let $t_{1}$ and $t_{2}$ be any numbers $t_{1}<f(x)<t_{2}$. With $I=\left(t_{1}, t_{2}\right)$ we have $D\left(f^{-1}(I), x\right)=1$. Hence $D\left[\mathbb{R}^{n}-f^{-1}(I), x\right]=0$. This implies $D\left(A_{t_{2}}, x\right)=D\left(B_{t_{1}}, x\right)=0$, which implies that the approximate limit of $f$ at $x$ is $f(x)$.

The next result shows the great similarity between continuity and approximate continuity.
7.54. Theorem. $f$ is approximately continuous at $x$ if and only if there exists a Lebesgue measurable set $E$ containing $x$ such that $D(E, x)=1$ and the restriction of $f$ to $E$ is continuous at $x$.

Proof. We will prove only the difficult direction. The other direction is left to the reader. Thus, assume that $f$ is approximately continuous at $x$. The definition of approximate continuity implies that there are positive numbers $r_{1}>r_{2}>r_{3}>\ldots$ tending to zero such that

$$
\lambda\left[B(x, r) \cap\left\{y:|f(y)-f(x)|>\frac{1}{k}\right\}\right]<\frac{\lambda[B(x, r)]}{2^{k}}, \quad \text { for } r \leq r_{k}
$$

Define

$$
E=\mathbb{R}^{n} \backslash \bigcup_{k=1}^{\infty}\left[\left\{B\left(x, r_{k}\right) \backslash B\left(x, r_{k+1}\right)\right\} \cap\left\{y:|f(y)-f(x)|>\frac{1}{k}\right\}\right] .
$$

From the definition of $E$, it follows that the restriction of $f$ to $E$ is continuous at $x$. In order to complete the assertion, we will show that $D(\widetilde{E}, x)=0$. For this purpose, choose $\varepsilon>0$ and let $J$ be such that $\sum_{k=J}^{\infty} \frac{1}{2^{k}}<\varepsilon$. Furthermore, choose $r$ such that $0<r<r_{J}$ and let $K \geq J$ be that integer such that $r_{K+1} \leq r<r_{K}$. Then,

$$
\begin{aligned}
\lambda\left[\left(\mathbb{R}^{n} \backslash E\right) \cap B(x, r)\right] \leq & \lambda\left[B(x, r) \cap\left\{y:|f(y)-f(x)|>\frac{1}{K}\right\}\right] \\
& +\sum_{k=K+1}^{\infty} \lambda\left[\left\{B\left(x, r_{k}\right) \backslash B\left(x, r_{k+1}\right)\right\}\right. \\
& \left.\cap\left\{y:|f(y)-f(x)|>\frac{1}{k}\right\}\right] \\
\leq & \frac{\lambda[B(x, r)]}{2^{K}}+\sum_{k=K+1}^{\infty} \frac{\lambda\left[B\left(x, r_{k}\right)\right]}{2^{k}} \\
\leq & \frac{\lambda[B(x, r)]}{2^{K}}+\sum_{k=K+1}^{\infty} \frac{\lambda[B(x, r)]}{2^{k}} \\
\leq & \lambda[B(x, r)] \sum_{k=K}^{\infty} \frac{1}{2^{k}} \\
\leq & \lambda[B(x, r)] \cdot \varepsilon,
\end{aligned}
$$

which yields the desired result since $\varepsilon$ is arbitrary.
7.55. Theorem. Assume $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is Lebesgue measurable. Then $f$ is approximately continuous $\lambda$ - almost everywhere.

Proof. First, we will prove that there exist disjoint compact sets $K_{i} \subset \mathbb{R}^{n}$ such that

$$
\lambda\left[\mathbb{R}^{n}-\bigcup_{i=1}^{\infty} K_{i}\right]=0
$$

and $f$ restricted to each $K_{i}$ is continuous. To this end, set $B_{i}=B(0, i)$ for each positive integer $i$. By Lusin's Theorem, there exists a compact set $K_{1} \subset B_{1}$ with $\lambda\left(B_{1}-K_{1}\right) \leq 1$ such that $f$ restricted to $K_{1}$ is continuous. Assuming that $K_{1}, K_{2}, \ldots, K_{j}$ have been constructed, appeal to Lusin's Theorem again to obtain a compact set $K_{j+1}$ such that

$$
K_{j+1} \subset B_{j+1}-\bigcup_{i=1}^{j} K_{i}, \quad \lambda\left[B_{j+1}-\bigcup_{i=1}^{j+1} K_{i}\right] \leq \frac{1}{j+1},
$$

and $f$ restricted to $K_{j+1}$ is continuous. Let

$$
E_{i}=K_{i} \cap\left\{x: D\left(K_{i}, x\right)=1\right\}
$$

and recall from Theorem 7.13 that $\lambda\left(K_{i}-E_{i}\right)=0$. Thus, $E_{i}$ has the property that $D\left(E_{i}, x\right)=1$ for each $x \in E_{i}$. Furthermore, $f$ restricted to $E_{i}$ is continuous since
$E_{i} \subset K_{i}$. Hence, by Theorem 7.54 we have that $f$ is approximately continuous at each point of $E_{i}$ and therefore at each point of

$$
\bigcup_{i=1}^{\infty} E_{i} .
$$

Since

$$
\lambda\left[\mathbb{R}^{n}-\bigcup_{i=1}^{\infty} E_{i}\right]=\lambda\left[\mathbb{R}^{n}-\bigcup_{i=1}^{\infty} K_{i}\right]=0
$$

we obtain the conclusion of the theorem.

## Exercises for Section 7.9

1. Define a set $E$ to be "density open" if $E$ is Lebesgue measurable and if $D(E, x)=$ 1 for all $x \in E$. Prove that the density open sets form a topology. The issue here is the following: In order to show that the "density open" sets form a topology, let $\left\{E_{\alpha}\right\}$ denote an arbitrary (possibly uncountable) collection of density open sets. It must be shown that $E:=\cup_{\alpha} E_{\alpha}$ is density open. In particular, it must be shown that $E$ is measurable.
2. Prove that a function $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is approximately continuous at every point if and only if $f$ is continuous in the density open topology of Exercise 1, Section 7.9 .
3. Let $f$ be an arbitrary function with the property that for each $x \in \mathbb{R}^{n}$ there exists a measurable set $E$ such that $D(E, x)=1$ and $f$ restricted to $E$ is continuous at $x$. Prove that $f$ is a measurable function
4. Suppose $f$ is a bounded measurable function on $\mathbb{R}^{n}$ that is approximately continuous at $x_{0}$. Prove that $x_{0}$ is a Lebesgue point for $f$. Hint: Use Definition 7.51 .
5. Show that if $f$ has a Lebesgue point at $x_{0}$, then $f$ is approximately continuous at $x_{0}$.

## CHAPTER 8

## Elements of Functional Analysis

### 8.1. Normed Linear Spaces

We have already encountered examples of normed linear spaces, namely the $L^{p}$ spaces. Here we introduce the notion of abstract normed linear spaces and begin the investigation of the structure of such spaces.
8.1. Definition. A linear space (or vector space) is a set $X$ that is endowed with two operations, addition and scalar multiplication, that satisfy the following conditions: for every $x, y, z \in X$ and $\alpha, \beta \in \mathbb{R}$
(i) $x+y=y+x \in X$.
(ii) $x+(y+z)=(x+y)+z$.
(iii) There is an element $0 \in X$ such that $x+0=x$ for each $x \in X$.
(iv) For each $x \in X$ there is an element $w \in X$ such that $x+w=0$.
(v) $\alpha x \in X$.
(vi) $\alpha(\beta x)=(\alpha \beta) x$.
(vii) $\alpha(x+y)=\alpha x+\alpha y$.
(viii) $(\alpha+\beta) x=\alpha x+\beta y$.
(ix) $1 x=x$.

We note here some immediate consequences of the definition of a linear space. If $x, y, z \in X$ and

$$
x+y=x+z
$$

then, by conditions (i) and (iv), there is a $w \in X$ such that $w+x=0$ and hence

$$
y=0+y=(w+x)+y=w+(x+y)=w+(x+z)=(w+x)+z=0+z=z
$$

Thus, in particular, for each $x \in X$ there is exactly one element $w \in X$ such that $x+w=0$. We will denote that element by $-x$ and we will write $y-x$ for $y+(-x)$.

If $\alpha \in \mathbb{R}$ and $x \in X$, then $\alpha x=\alpha(x+0)=\alpha x+\alpha 0$, from which we can conclude $\alpha 0=0$. Similarly $\alpha x=(\alpha+0) x=\alpha x+0 x$, from which we conclude $0 x=0$.

If $\lambda \neq 0$ and $\lambda x=0$, then

$$
x=\left(\frac{\lambda}{\lambda}\right) x=\frac{1}{\lambda}(\lambda x)=\frac{1}{\lambda} 0=0
$$

8.2. Definition. A subset $Y$ of a linear space $X$ is a subspace of $X$ if $\alpha x+\beta y \in$ $Y$ for all $x, y \in Y$ and $\alpha, \beta \in \mathbb{R}$.

Thus if $Y$ is a subspace of $X$, then $Y$ is itself a linear space with respect to the addition and scalar multiplication it inherits from $X$. The notion of subspace that we have defined above might more properly be called linear subspace to distinguish it from the notion of topological subspace in case $X$ is also a topological space. In this chapter we will use the term "subspace" in the sense of the definition above. If we have occasion to refer to a topological subspace we will mention it explicitly.

If $S$ is a nonempty subspace of a linear space $X$, the set $Y$ of all elements of $X$ of the form

$$
\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{m} x_{m}
$$

where $m$ is any positive integer and $\alpha_{j} \in \mathbb{R}, x_{j} \in S$ for $1 \leq j \leq m$ is easily seen to be a subspace of $X$. The subspace $Y$ will be called the subspace spanned by $S$. It is the smallest subspace of $X$ that contains $S$.
8.3. Definition. If $S=\left\{x_{1}, x_{2}, \ldots, x_{m}\right\}$ is a finite subset of a linear space $X$, then $S$ is linearly independent if

$$
\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{m} x_{m}=0
$$

implies $\alpha_{1}=\alpha_{2}=\cdots=\alpha_{m}=0$. In general, a subset $S$ of $X$ is linearly independent if every finite subset of $S$ is linearly independent.

Suppose $S$ is a subset of a linear space $X$. If $S$ is linearly independent and $X$ is spanned by $S$, then for any $x \in X$ there is a finite subset $\left\{x_{i}\right\}_{i=1}^{m}$ of $S$ and a finite sequence of real numbers $\left\{\alpha_{i}\right\}_{i=1}^{m}$ such that

$$
x=\sum_{i=1}^{m} \alpha_{i} x_{i}
$$

where $\alpha_{i} \neq 0$ for each $i$.
Suppose that for some other choice $\left\{y_{j}\right\}_{j=1}^{k}$ of elements in $S$ and real numbers $\left\{\beta_{j}\right\}_{j=1}^{k}$

$$
x=\sum_{j=1}^{k} \beta_{j} y_{j}
$$

where $\beta_{j} \neq 0$ for each $j$. Then

$$
\sum_{i=1}^{m} \alpha_{i} x_{i}-\sum_{j=1}^{k} \beta_{j} y_{j}=0
$$

If for some $i \in\{1,2, \ldots, m\}$ there were no $j \in\{1.2 \ldots, k\}$ such that $x_{i}=y_{j}$, then since $S$ is linearly independent we would have $\alpha_{i}=0$. Similarly, for each $j \in$
$\{1.2 \ldots, k\}$ there is an $i \in\{1,2, \ldots, m\}$ such that $x_{i}=y_{j}$. Thus the two sequences $\left\{x_{i}\right\}_{i=1}^{m}$ and $\left\{y_{j}\right\}_{j=1}^{k}$ must contain exactly the same elements. Renumbering the $\beta_{j}$, if necessary, we see that

$$
\sum_{i=1}^{m}\left(\alpha_{i}-\beta_{i}\right) x_{i}=0
$$

Since $S$ is linearly independent we have $\alpha_{i}=\beta_{i}$ for $1 \leq i \leq k$. Thus each $x \in X$ has a unique representation as a finite linear combination of elements of $S$.
8.4. Definition. A subset of a linear space $X$ that is linearly independent and spans $X$ is a basis for $X$. A linear space is finite dimensional if it has a finite basis.

The proof of the following result is a consequence of the Hausdorff Maximal Principle; see Exercise 1, Section 8.1.

### 8.5. Theorem. Every linear space has a basis.

8.6. Examples. (i) The set $\mathbb{R}^{n}$ is a linear space with respect to the addition and scalar multiplication defined by

$$
\begin{aligned}
\left(x_{1}, x_{2}, \ldots, x_{n}\right)+\left(y_{1}, y_{2}, \ldots, y_{n}\right) & =\left(x_{1}+y_{1}, \ldots, x_{n}+y_{n}\right) \\
\alpha\left(x_{1}, x_{2}, \ldots, x_{n}\right) & =\left(\alpha x_{1}, \alpha x_{2}, \ldots, \alpha x_{n}\right)
\end{aligned}
$$

(ii) If $A$ is any set, the set of all real-valued functions on $A$ is a linear space with respect to the addition and scalar multiplication

$$
\begin{aligned}
(f+g)(x) & =f(x)+g(x) \\
(\alpha f)(x) & =\alpha f(x)
\end{aligned}
$$

(iii) If $S$ is a topological space, then the set $C(S)$ of all real-valued continuous functions on $S$ is a linear space with respect to the addition and scalar multiplication defined in (ii).
(iv) If $(X, \mu)$ is a measure space, then by, Lemma Theorem $6.23, L^{p}(X, \mu)$ is a linear space for $1 \leq p \leq \infty$ with respect to the addition and scalar multiplication defined in (ii).
(v) For $1 \leq p<\infty$ let $l^{p}$ denote the set of all sequences $\left\{a_{k}\right\}_{k=1}^{\infty}$ of real numbers such that $\sum_{k=1}^{\infty}\left|a_{k}\right|^{p}$ converges. Each such sequence may be viewed as a real-valued function on the set of positive integers. If addition and scalar are defined as in example (ii), then each $l^{p}$ is a linear space.
(vi) Let $l^{\infty}$ denote the set of all bounded sequences of real numbers. Then with respect the addition and scalar multiplication of $(\mathrm{v}), l^{\infty}$ is a linear space.
8.7. Definition. Let $X$ be a linear space. A function $\|\cdot\|: X \rightarrow \mathbb{R}$ is a norm on $X$ if
(i) $\|x+y\| \leq\|x\|+\|y\|$ for all $x, y \in X$,
(ii) $\|\alpha x\|=|\alpha|\|x\|$ for all $x \in X$ and $\alpha \in \mathbb{R}$,
(iii) $\|x\| \geq 0$ for each $x \in X$,
(iv) $\|x\|=0$ only if $x=0$.

A real-valued function on $X$ satisfying conditions (i), (ii), and (iii) is a semi-norm on $X$. A linear space $X$ equipped with a norm $\|\cdot\|$ is a normed linear space.

Suppose $X$ is a normed linear space with norm $\|\cdot\|$. For $x, y \in X$ set

$$
\rho(x, y)=\|x-y\|
$$

Then $\rho$ is a nonnegative real-valued function on $X \times X$, and from the properties of $\|\cdot\|$ we see that
(i) $\rho(x, y)=0$ if and only if $x=y$,
(ii) $\rho(x, y)=\rho(y, x)$ for any $x, y \in X$,
(iii) $\rho(x, z) \leq \rho(x, y)+\rho(y, z)$ for any $x, y, z \in X$.

Thus $\rho$ is a metric on $X$, and we see that a normed linear space is also a metric space; in particular, it is a topological space. Functional Analysis (in normed linear spaces) is essentially the study of the interaction between the algebraic (linear) structure and the topological (metric) structure of such spaces.

In a normed linear space $X$ we will denote by $B(x, r)$ the open ball with center at $x$ and radius $r$, i.e.,

$$
B(x, r)=\{y \in X:\|y-x\|<r\}
$$

8.8. Definition. A normed linear space is a Banach space if it is a complete metric space with respect to the metric induced by its norm.
8.9. Examples. (i) $\mathbb{R}^{n}$ is a Banach space with respect to the norm

$$
\left\|\left(x_{1}, x_{2}, \ldots, x_{n}\right)\right\|=\left(x_{1}^{2}+x_{2}^{2}+\cdots+x_{n}^{2}\right)^{\frac{1}{2}}
$$

(ii) For $1 \leq p \leq \infty$ the linear spaces $L^{p}(X, \mu)$ are Banach spaces with respect to the norms

$$
\begin{aligned}
\|f\|_{L^{p}(X, \mu)} & =\left(\int_{X}|f|^{p} d \mu\right)^{\frac{1}{p}} \quad \text { if } \quad 1 \leq p<\infty \\
\|f\|_{L^{p}(X, \mu)} & =\inf \{M: \mu(\{x \in X:|f(x)|>M\})=0\} \quad \text { if } \quad p=\infty
\end{aligned}
$$

This is a rephrasing of Theorem 6.29.
(iii) For $1 \leq p \leq \infty$ the linear spaces $l^{p}$ of Examples 8.6(v), (vi) are Banach spaces with respect to the norms

$$
\begin{aligned}
& \left\|\left\{a_{k}\right\}\right\|_{l^{p}}=\left(\sum_{k=1}^{\infty}\left|a_{k}\right|^{p}\right)^{\frac{1}{p}} \quad \text { if } \quad 1 \leq p<\infty \\
& \left\|\left\{a_{k}\right\}\right\|_{l^{p}}=\sup _{k \geq 1}\left|a_{k}\right| \quad \text { if } \quad p=\infty
\end{aligned}
$$

This is a consequence of Theorem 6.29 with appropriate choices of $X$ and $\mu$. (iv) If $X$ is a compact metric space, then the linear space $C(X)$ of all continuous real-valued functions on $X$ is a Banach space with respect to the norm

$$
\|f\|=\sup _{x \in X}|f(x)|
$$

If $\left\{x_{k}\right\}$ is a sequence in a Banach space $X$, the series $\sum_{k=1}^{\infty} x_{k}$ converges to $x \in X$ if the sequence of partial sums $s_{m}=\sum_{k=1}^{m} x_{k}$ converges to $x$, i.e.,

$$
\left\|x-s_{m}\right\|=\left\|x-\sum_{k=1}^{m} x_{k}\right\| \rightarrow 0
$$

as $m \rightarrow \infty$.
8.10. Proposition. Suppose $X$ is a Banach space. Then any absolutely convergence series is convergent, i.e., if the series $\sum_{k=1}^{\infty}\left\|x_{k}\right\|$ converges in $\mathbb{R}$, then the series $\sum_{k=1}^{\infty} x_{k}$ converges in $X$.

Proof. If $m>l$, then

$$
\left\|\sum_{k=1}^{m} x_{k}-\sum_{k=1}^{l} x_{k}\right\|=\left\|\sum_{k=l+1}^{m} x_{k}\right\| \leq \sum_{k=l+1}^{m}\left\|x_{k}\right\| .
$$

Thus if $\sum_{k=1}^{\infty}\left\|x_{k}\right\|$ converges in $\mathbb{R}$ the sequence of partial sums $\left\{\sum_{k=1}^{m} x_{k}\right\}_{m=1}^{\infty}$ is Cauchy in $X$ and therefore converges to some element.
8.11. Definitions. Suppose $X$ and $Y$ are linear spaces. A mapping $T: X \rightarrow Y$ is linear if for every $x, y \in X$ and $\alpha, \beta \in \mathbb{R}$

$$
T(\alpha x+\beta y)=\alpha T(x)+\beta T(y)
$$

If $X$ and $Y$ are normed linear spaces and $T: X \rightarrow Y$ is a linear mapping, then $T$ is bounded if there exists a constant $M$ such that

$$
\|T(x)\| \leq M\|x\|
$$

for each $x \in X$.
8.12. Theorem. Suppose $X$ and $Y$ are normed linear spaces and $T: X \rightarrow Y$ is linear. Then
(i) The linear mapping $T$ is bounded if and only if

$$
\sup \{\|T(x)\|: x \in X,\|x\|=1\}<\infty
$$

(ii) The linear mapping $T$ is continuous if and only if it is bounded.

Proof. (i) If $M$ is a constant such that

$$
\|T(x)\| \leq M\|x\|
$$

for each $x \in X$, then for any $x \in X$ with $\|x\|=1$ we have $\|T(x)\| \leq M$.
On the other hand, if

$$
K=\sup \{\|T(x)\|: x \in X,\|x\|=1\}<\infty
$$

then for any $0 \neq x \in X$ we have

$$
\|T(x)\|=\|x\|\left\|T\left(\frac{x}{\|x\|}\right)\right\| \leq K\|x\|
$$

(ii) If $T$ is continuous at 0 , then there exists a $\delta>0$ such that $\|T(x)\| \leq 1$ whenever $x \in B(0,2 \delta)$. If $\|x\|=1$, then

$$
\|T(x)\|=\frac{1}{\delta}\|T(\delta x)\| \leq \frac{1}{\delta}
$$

Thus

$$
\sup \{\|T(x)\|: x \in X,\|x\|=1\} \leq \frac{1}{\delta}
$$

On the other hand, if there is a constant $M$ such that

$$
\|T(x)\| \leq M\|x\|
$$

for each $x \in X$, then for any $\varepsilon>0$

$$
\|T(x)\|<\varepsilon
$$

whenever $x \in B\left(0, \frac{\varepsilon}{M}\right)$. Let $x_{0} \in X$. If $\left\|y-x_{0}\right\|<\frac{\varepsilon}{M}$, then

$$
\left\|T(y)-T\left(x_{0}\right)\right\|=\left\|T\left(y-x_{0}\right)\right\|<\varepsilon
$$

Thus $T$ is continuous on $X$.
8.13. Definition. If $T: X \rightarrow Y$ is a linear mapping from a normed linear space $X$ into a normed linear space $Y$ we set

$$
\begin{equation*}
\|T\|=\sup \{\|T(x)\|: x \in X,\|x\|=1\} \tag{8.1}
\end{equation*}
$$

This choice of notation will be justified by Theorem 8.16, where we will show that $\|T\|$ is a norm on an appropriate linear space.
8.14. Proposition. Suppose $X$ and $Y$ are normed linear spaces, $\left\{T_{k}\right\}$ is a sequence of bounded linear mappings of $X$ into $Y$, and $T: X \rightarrow Y$ is a mapping such that

$$
\lim _{k \rightarrow \infty}\left\|T_{k}(x)-T(x)\right\|=0
$$

for each $x \in X$. Then $T$ is a linear mapping and

$$
\|T\| \leq \liminf _{k \rightarrow \infty}\left\|T_{k}\right\|
$$

Proof. For $x, y \in X$ and $\alpha, \beta \in \mathbb{R}$

$$
\begin{aligned}
\|T(\alpha x+\beta y)-\alpha T(x)-\beta T(y)\| \leq & \left\|T(\alpha x+\beta y)-T_{k}(\alpha x+\beta y)\right\| \\
& +\left\|\alpha\left(T_{k}(x)-T(x)\right)+\beta\left(T_{k}(y)-T(y)\right)\right\| \\
\leq & \left\|T(\alpha x+\beta y)-T_{k}(\alpha x+\beta y)\right\| \\
& +|\alpha|\left\|T_{k}(x)-T(x)\right\|+|\beta|\left\|T_{k}(y)-T(y)\right\|
\end{aligned}
$$

for any $k \geq 1$. Thus $T$ is a linear mapping.
If $x \in X$ with $\|x\|=1$, then

$$
\|T(x)\| \leq\left\|T_{k}(x)\right\|+\left\|T(x)-T_{k}(x)\right\| \leq\left\|T_{k}\right\|+\left\|T(x)-T_{k}(x)\right\|
$$

for any $k \geq 1$. Thus

$$
\|T(x)\| \leq \liminf _{k \rightarrow \infty}\left\|T_{k}\right\|
$$

for any $x \in X$ with $\|x\|=1$ and hence

$$
\|T\|=\sup \{\|T(x)\|: x \in X,\|x\|=1\} \leq \liminf _{k \rightarrow \infty}\left\|T_{k}\right\|
$$

Suppose $X$ and $Y$ are linear spaces, and let $\mathcal{L}(X, Y)$ denote the set of all linear mappings of $X$ into $Y$. For $T, S \in \mathcal{L}(X, Y)$ and $\alpha, \beta \in \mathbb{R}$ define

$$
\begin{aligned}
(T+S)(x) & =T(x)+S(x) \\
(\alpha T)(x) & =\alpha T(x)
\end{aligned}
$$

for $x \in X$. Note that these are the "usual" definitions of the sum and scalar multiple of functions. It is left as an exercise to show that with these operations $\mathcal{L}(X, Y)$ is a linear space.
8.15. Notation. If $X$ and $Y$ are normed linear spaces, denote by $\mathcal{B}(X, Y)$ the set of all bounded linear mappings of $X$ into $Y$. Clearly $\mathcal{B}(X, Y)$ is a subspace of $\mathcal{L}(X, Y)$. In case $Y=\mathbb{R}$ we will refer to the elements of $\mathcal{L}(X, \mathbb{R})$ as linear functionals on $X$.
8.16. Theorem. Suppose $X$ and $Y$ are normed linear spaces. Then (8.1) defines a norm on $\mathcal{B}(X, Y)$. If $Y$ is a Banach space, then $\mathcal{B}(X, Y)$ is a Banach space with respect to this norm.

Proof. Clearly $\|T\| \geq 0$. If $\|T\|=0$, then $T(x)=0$ for any $x \in X$ with $\|x\|=1$. Thus if $0 \neq x \in X$, then

$$
T(x)=\|x\| T\left(\frac{x}{\|x\|}\right)=0
$$

i.e., $T=0$.

If $T, S \in \mathcal{B}(X, Y)$ and $\alpha \in \mathbb{R}$, then

$$
\|\alpha T\|=\sup \{|\alpha|\|T(x)\|: x \in X,\|x\|=1\}=|\alpha|\|T\|
$$

and

$$
\begin{aligned}
\|T+S\| & =\sup \{\|T(x)+S(x)\|: x \in X,\|x\|=1\} \\
& \leq \sup \{\|T(x)\|+\|S(x)\|: x \in X,\|x\|=1\} \\
& \leq\|T\|+\|S\|
\end{aligned}
$$

Thus (8.1) defines a norm on $\mathcal{B}(X, Y)$.
Suppose $Y$ is a Banach space and $\left\{T_{k}\right\}$ is a Cauchy sequence in $\mathcal{B}(X, Y)$. Then $\left\{\left\|T_{k}\right\|\right\}$ is bounded, i.e., there is a constant $M$ such that $\left\|T_{k}\right\| \leq M$ for all $k$.

For any $x \in X$ and $k, m \geq 1$

$$
\left\|T_{k}(x)-T_{m}(x)\right\| \leq\left\|T_{k}-T_{m}\right\|\|x\|
$$

Thus $\left\{T_{k}(x)\right\}$ is a Cauchy sequence in $Y$. Since $Y$ is a Banach space there is an element $T(x) \in Y$ such that

$$
\left\|T_{k}(x)-T(x)\right\| \rightarrow 0
$$

as $k \rightarrow \infty$. In view of Proposition 8.14 we know that $T$ is a linear mapping of $X$ into $Y$. Moreover, again by Proposition 8.14

$$
\|T\|=\liminf _{k \rightarrow \infty}\left\|T_{k}\right\| \leq M
$$

whence $T \in \mathcal{B}(X, Y)$.

## Exercises for Section 8.1

1. Use the Hausdorff Maximal Principle to show that every linear space has a basis. Hint: observe that a linearly independent subset of a linear space $X$ spans $X$ if, and only if, it is maximal with respect to set inclusion i.e., it is not contained in any other linearly independent subset of $X$.
2. For $i=1,2, \ldots, m$, let $X_{i}$ be a Banach space with norm $\|\cdot\|_{i}$. The Cartesian product

$$
X=\prod_{i=1}^{m} X_{i}
$$

consisting of points $x=\left(x_{1}, x_{2}, \ldots, x_{m}\right)$ with $x_{i} \in X_{i}$ is a vector space under the definitions

$$
x+y=\left(x_{1}+y_{1}, \ldots, x_{m}+y_{m}\right), \quad c x=\left(c x_{1}, \ldots, c x_{m}\right) .
$$

Prove that $X$ is a Banach space with respect to any of the equivalent norms

$$
\|x\|=\left(\sum_{i=1}^{m}\left\|x_{i}\right\|_{i}^{p}\right)^{1 / p}, \quad 1 \leq p<\infty
$$

### 8.2. Hahn-Banach Theorem

In this section we prove the existence of extensions of linear functionals from a subspace $Y$ of $X$ to all of $X$ satisfying various conditions.
8.17. Theorem (Hahn-Banach Theorem: semi-norm version). Suppose $X$ is a linear space and $p$ is a semi-norm on $X$. Let $Y$ be a subspace of $X$ and $f: Y \rightarrow \mathbb{R}$ a linear functional such that $f(x) \leq p(x)$ for all $x \in Y$. Then there exists a linear functional $g: X \rightarrow \mathbb{R}$ such that $g(x)=f(x)$ for all $x \in Y$ and $g(x) \leq p(x)$ for all $x \in X$.

Proof. Let $\mathcal{F}$ denote the family of all pairs $(W, h)$ where $W$ is a subspace with $Y \subset W \subset X$ and $h$ is a linear functional on $W$ such that $h=f$ on $Y$ and $h \leq p$ on $W$. For each $(W, h) \in \mathcal{F}$ let

$$
G(W, h)=\{(x, r): x \in W, r=h(x)\} \subset X \times \mathbb{R}
$$

Observe that if $\left(W_{1}, h_{1}\right),\left(W_{2}, h_{2}\right) \in \mathcal{F}$, then $G\left(W_{1}, h_{1}\right) \subset G\left(W_{2}, h_{2}\right)$ if and only if

$$
\begin{align*}
W_{1} & \subset W_{2}  \tag{8.2}\\
h_{2} & =h_{1} \text { on } W_{1} . \tag{8.3}
\end{align*}
$$

Set $\mathcal{E}=\{G(W, h):(W, h) \in \mathcal{F}\}$. If $\mathcal{T}$ is a subfamily of $\mathcal{E}$ that is linearly ordered by inclusion, set

$$
W_{\infty}=\bigcup\{W: G(W, h) \in \mathcal{T} \text { for some }(W, h) \in \mathcal{T}\}
$$

Clearly $W_{\infty}$ is a subspace with $Y \subset W_{\infty} \subset X$. In view of (8.2) we can define a linear functional $h_{\infty}$ on $W_{\infty}$ by setting

$$
h_{\infty}(x)=h(x) \quad \text { if } \quad G(W, h) \in \mathcal{T} \text { and } x \in W
$$

Thus $\left(W_{\infty}, h_{\infty}\right) \in \mathcal{F}$ and $G(W, h) \subset G\left(W_{\infty}, h_{\infty}\right)$ for each $G(W, h) \in \mathcal{T}$. Hence, we may apply Zorn's Lemma (see p. 8), to conclude that $\mathcal{E}$ contains a maximal element. This means there is a pair $\left(W_{0}, h_{0}\right) \in \mathcal{F}$ such that $G\left(W_{0}, h_{0}\right)$ is not contained in any other set $G(W, h) \in \mathcal{E}$.

Since from the definition of $\mathcal{F}$ we have $h_{0}=f$ on $Y$ and $h_{0} \leq p$ on $W_{0}$, the proof will be complete when we show that $W_{0}=X$.

To the contrary, suppose $W_{0} \neq X$, and let $x_{0} \in X-W_{0}$. Set

$$
W^{\prime}=\left\{x+\alpha x_{0}: x \in W_{0}, \alpha \in \mathbb{R}\right\}
$$

Then $W^{\prime}$ is a subspace of $X$ containing $W_{0}$. If $x, x^{\prime} \in W_{0}, \alpha, \alpha^{\prime} \in \mathbb{R}$ and

$$
x+\alpha x_{0}=x^{\prime}+\alpha^{\prime} x_{0}
$$

then

$$
x-x^{\prime}=\left(\alpha^{\prime}-\alpha\right) x_{0}
$$

If $\alpha^{\prime}-\alpha \neq 0$ this would imply that $x_{0} \in W_{0}$. Thus $x=x^{\prime}$ and $\alpha^{\prime}=\alpha$. Thus each element of $W^{\prime}$ has a unique representation in the form $x+\alpha x_{0}$, and hence we can define a linear functional $h^{\prime}$ on $W^{\prime}$ by fixing $c \in \mathbb{R}$ and setting

$$
h^{\prime}\left(x+\alpha x_{0}\right)=h_{0}(x)+\alpha c
$$

for each $x \in W_{0}, \alpha \in \mathbb{R}$. Clearly $h^{\prime}=h_{0}$ on $W_{0}$. We now choose $c$ so that $h^{\prime} \leq p$ on $W^{\prime}$.

Observe that for $x, x^{\prime} \in W_{0}$

$$
h_{0}(x)+h_{0}\left(x^{\prime}\right)=h_{0}\left(x+x_{0}+x^{\prime}-x_{0}\right) \leq p\left(x+x_{0}\right)+p\left(x^{\prime}-x_{0}\right)
$$

and thus

$$
h_{0}\left(x^{\prime}\right)-p\left(x^{\prime}-x_{0}\right) \leq p\left(x+x_{0}\right)-h_{0}(x)
$$

for any $x, x^{\prime} \in W_{0}$.

In view of the last inequality there is a $c \in \mathbb{R}$ such that

$$
\sup \left\{h_{0}(x)-p\left(x-x_{0}\right): x \in W_{0}\right\} \leq c \leq \inf \left\{p\left(x+x_{0}\right)-h_{0}(x): x \in W_{0}\right\}
$$

With this choice of $c$ we see that for any $x \in W_{0}$ and $\alpha \neq 0$

$$
h^{\prime}\left(x+\alpha x_{0}\right)=\alpha h^{\prime}\left(\frac{x}{\alpha}+x_{0}\right)=\alpha\left(h_{0}\left(\frac{x}{\alpha}\right)+c\right)
$$

if $\alpha>0$, then $\alpha c \leq p\left(x+x_{0}\right)-h_{0}(x)$ and therefore

$$
\begin{aligned}
h^{\prime}\left(x+\alpha x_{0}\right) & \leq \alpha\left(h_{0}\left(\frac{x}{\alpha}\right)+p\left(\frac{x}{\alpha}+x_{0}\right)-h_{0}\left(\frac{x}{\alpha}\right)\right) \\
& =\alpha p\left(\frac{x}{\alpha}+x_{0}\right)=p\left(x+\alpha x_{0}\right) .
\end{aligned}
$$

If $\alpha<0$, then

$$
\begin{aligned}
h^{\prime}\left(x+\alpha x_{0}\right) & =\alpha\left(h_{0}\left(\frac{x}{\alpha}\right)+c\right) \\
& =|\alpha|\left(h_{0}\left(\frac{x}{|\alpha|}\right)-c\right) \\
& \leq|\alpha|\left(h_{0}\left(\frac{x}{|\alpha|}\right)+p\left(\frac{x}{|\alpha|}-x_{0}\right)-h_{0}\left(\frac{x}{|\alpha|}\right)\right) \\
& =|\alpha| p\left(\frac{x}{|\alpha|}-x_{0}\right)=p\left(x+\alpha x_{0}\right) .
\end{aligned}
$$

Thus $\left(W^{\prime}, h^{\prime}\right) \in \mathcal{F}$ and $G\left(W_{0}, h_{0}\right)$ is a proper subset of $G\left(W^{\prime}, h^{\prime}\right)$, contradicting the maximality of $G\left(W_{0}, h_{0}\right)$. This implies our assumption that $W_{0} \neq X$ must be false and so it follows that $g:=h_{0}$ is a linear functional on $X$ such that $g=f$ on $Y$ and $g \leq p$ on $X$.
8.18. Remark. The proof of Theorem 8.17 actually gives more than is asserted in the statement of the Theorem. A careful reading of the proof shows that the function $p$ need not be a semi-norm; it suffices that $p$ be subadditive, i.e.,

$$
p(x+y) \leq p(x)+p(y)
$$

and positively homogeneous, i.e.,

$$
p(\alpha x)=\alpha p(x)
$$

whenever $\alpha \geq 0$. In particular $p$ need not be nonnegative.
As an immediate consequence of Theorem 8.17 we obtain
8.19. Theorem (Hahn-Banach Theorem: norm version). Suppose $X$ is a normed linear space and $Y$ is a subspace of $X$. If $f$ is a linear functional on $Y$ and $M$ is a positive constant such that

$$
|f(x)| \leq M\|x\|
$$

for each $x \in Y$, then there is a linear functional $g$ on $X$ such that $g=f$ on $Y$ and

$$
|g(x)| \leq M\|x\|
$$

for each $x \in X$.
Proof. Observe that $p(x)=M\|x\|$ is a semi-norm on $X$. Thus by Theorem 8.17 there is a linear functional $g$ on $X$ that extends $f$ to $X$ and such that

$$
g(x) \leq M\|x\|
$$

for each $x \in X$. Since $g(-x)=-g(x)$ and $\|-x\|=\|x\|$ it follows immediately that

$$
|g(x)| \leq M\|x\|
$$

for each $x \in X$.
The following is a useful consequence of Theorem 8.19,
8.20. Theorem. Suppose $X$ is a normed linear space, $Y$ is a subspace of $X$, and $x_{0} \in X$ such that

$$
\rho=\inf _{y \in Y}\left\|x_{0}-y\right\|>0
$$

i.e., the distance from $x_{0}$ to $Y$ is positive. Then there is a bounded linear functional $f$ on $X$ such that

$$
f(y)=0 \text { for all } y \in Y, \quad f\left(x_{0}\right)=1 \quad \text { and } \quad\|f\|=\frac{1}{\rho}
$$

Proof. We will use the following observation throughout the proof, namely, that since $Y$ is a vector space, then

$$
\rho=\inf _{y \in Y}\left\|x_{0}-y\right\|=\inf _{y \in Y}\left\|x_{0}-(-y)\right\|=\inf _{y \in Y}\left\|x_{0}+y\right\|
$$

Set

$$
W=\left\{y+\alpha x_{0}: y \in Y, \alpha \in \mathbb{R}\right\}
$$

Then, as noted in the proof of Theorem 8.17, $W$ is a subspace of $X$ containing $Y$ and each element of $W$ has a unique representation in the form $y+\alpha x_{0}$ with $y \in Y$ and $\alpha \in \mathbb{R}$. Thus we can define a linear functional $g$ on $W$ by

$$
g\left(y+\alpha x_{0}\right)=\alpha
$$

If $\alpha \neq 0$ then

$$
\left\|y+\alpha x_{0}\right\|=|\alpha|\left\|\frac{y}{\alpha}+x_{0}\right\| \geq|\alpha| \rho
$$

whence

$$
\left|g\left(y+\alpha x_{0}\right)\right| \leq \frac{1}{\rho}\left\|y+\alpha x_{0}\right\|
$$

Thus $g$ is a bounded linear functional on $W$ such that $g(y)=0$ for $y \in Y,|g(w)| \leq$ $\frac{1}{\rho}\|w\|$ for $w \in W$ and $g\left(x_{0}\right)=1$. In view of Theorem 8.19 there is a bounded linear
functional $f$ on $X$ such that $f=g$ on $W$ and $\|f\| \leq \frac{1}{\rho}$. There is a sequence $\left\{y_{k}\right\}$ in $Y$ such that

$$
\lim _{k \rightarrow \infty}\left\|y_{k}+x_{0}\right\|=\rho
$$

Let $x_{k}=\frac{y_{k}+x_{0}}{\left\|y_{k}+x_{0}\right\|}$. Then $\left\|x_{k}\right\|=1$ and

$$
\left\|f\left(x_{k}\right)\right\|=\left\|g\left(x_{k}\right)\right\|=\frac{1}{\left\|y_{k}+x_{0}\right\|}
$$

for all $k$, whence $\|f\|=\frac{1}{\rho}$.

## Exercises for Section 8.2

1. Suppose $Y$ is a closed subspace of a normed linear space $X, Y \neq X$, and $\varepsilon>0$.

Show that there is an element $x \in X$ such that $\|x\|=1$ and

$$
\inf _{y \in Y}\|x-y\|>1-\varepsilon
$$

2. Let $f: X \rightarrow Y$ be a linear mapping of a normed linear space $X$ into a normed linear space $Y$. Show that $f$ is bounded if and only if $f$ is continuous at one point.
3. The kernel of a linear $f: X \rightarrow \mathbb{R}^{1}$ is the set $\{x: f(x)=0\}$. Prove that $f$ is bounded if and only if the kernel of $f$ is closed in $X$.

### 8.3. Continuous Linear Mappings


#### Abstract

In this section we deduce from the Baire Category Theorem three important results concerning continuous linear mappings between Banach spaces.


We first prove a "linear" version of the Uniform Boundedness Principle; Theorem 3.35.
8.21. Theorem (Uniform Boundedness Principle). Let $\mathcal{F}$ be a family of continuous linear mappings from a Banach space $X$ into a normed linear space $Y$ such that

$$
\sup _{T \in \mathcal{F}}\|T(x)\|<\infty
$$

for each $x \in X$. Then

$$
\sup _{T \in \mathcal{F}}\|T\|<\infty
$$

Proof. Observe that for each $T \in \mathcal{F}$ the real-valued function $x \mapsto\|T(x)\|$ is continuous. In view of Theorem 3.35, there is a nonempty open subset $U$ of $X$ and an $M>0$ such that

$$
\|T(x)\| \leq M
$$

for each $x \in U$ and $T \in \mathcal{F}$.
Fix $x_{0} \in U$ and let $r>0$ be such that $B\left(x_{0}, r\right) \subset U$. If $z=x_{0}+y \in B\left(x_{0}, r\right)$ and $T \in \mathcal{F}$, then

$$
\|T(z)\|=\left\|T\left(x_{0}\right)+T(y)\right\| \geq\|T(y)\|-\left\|T\left(x_{0}\right)\right\|
$$

i.e.,

$$
\|T(y)\| \leq\|T(z)\|+\left\|T\left(x_{0}\right)\right\| \leq 2 M
$$

for each $y \in B(0, r)$.
If $x \in X$ with $\|x\|=1$, then $\rho x \in B(0, r)$ for all $0<\rho<r$, in particular for $\rho=r / 2$. Therefore, for each $T \in \mathcal{F}$

$$
\|T(x)\|=\frac{2}{r}\left\|T\left(\frac{r}{2} x\right)\right\| \leq \frac{4 M}{r}
$$

Thus

$$
\sup _{T \in \mathcal{F}}\|T\| \leq \frac{4 M}{r}
$$

8.22. Corollary. Suppose $X$ is a Banach space, $Y$ is a normed linear space, $\left\{T_{k}\right\}$ is a sequence of bounded linear mappings of $X$ into $Y$ and $T: X \rightarrow Y$ is a mapping such that

$$
\lim _{k \rightarrow \infty}\left\|T_{k}(x)-T(x)\right\|=0
$$

for each $x \in X$. Then $T$ is a bounded linear mapping and

$$
\|T\| \leq \liminf _{k \rightarrow \infty}\left\|T_{k}\right\|<\infty
$$

Proof. In view of Proposition 8.14 we only need to show that $\left\{\left\|T_{k}\right\|\right\}$ is bounded. Since for each $x \in X$ the sequence $\left\{T_{k}(x)\right\}$ converges, it is bounded. From Theorem 8.21 we see that $\left\{\left\|T_{k}\right\|\right\}$ is bounded.
8.23. Definition. Let $X$ and $Y$ be normed linear spaces. A mapping $T: X \rightarrow$ $Y$ is said to be open if $T(U)$ is an open subset of $Y$ whenever $U$ is an open subset of $X$.
8.24. Theorem (Open Mapping Theorem). If $T$ is a bounded linear mapping of a Banach space $X$ onto a Banach space $Y$, then $T$ is an open mapping.

Proof. Fix $\varepsilon>0$. Since $T$ maps $X$ onto $Y$

$$
\begin{aligned}
Y & =\bigcup_{k=1}^{\infty} T(B(0, k \varepsilon)) . \\
& =\bigcup_{k=1}^{\infty} \overline{T(B(0, k \varepsilon)) .} .
\end{aligned}
$$

Since $Y$ is a complete metric space, the Baire Category Theorem asserts that one of these closed sets has a non-empty interior; that is, there exist $k_{0} \geq 1, y \in Y$ and $\delta>0$ such that

$$
\overline{T\left(B\left(0, k_{0} \varepsilon\right)\right.} \supset B(y, \delta)
$$

First, we will show that the origin is in the interior of $\overline{T(B(0,2 \varepsilon))}$. For this purpose, note that if $z \in B\left(\frac{y}{k_{0}}, \frac{\delta}{k_{0}}\right)$ then

$$
\left\|z-\frac{y}{k_{0}}\right\|<\frac{\delta}{k_{0}}
$$

i.e.,

$$
\left\|k_{0} z-y\right\|<\delta
$$

Thus $k_{0} z \in B(y, \delta) \subset \overline{T\left(B\left(0, k_{0} \varepsilon\right)\right)}$ which implies that $z \in \overline{T(B(0, \varepsilon))}$. Setting $y_{0}=\frac{y}{k_{0}}$ and $\delta_{0}=\frac{\delta}{k_{0}}$ we have

$$
B\left(y_{0}, \delta_{0}\right) \subset \overline{T(B(0, \varepsilon))}
$$

If $w \in B\left(0, \delta_{0}\right)$, then $z=y_{0}+w \in B\left(y_{0}, \delta_{0}\right)$ and there exist sequences $\left\{x_{k}\right\}$ and $\left\{x_{k}^{\prime}\right\}$ in $B(0, \varepsilon)$ such that

$$
\begin{aligned}
\left\|T\left(x_{k}\right)-y_{0}\right\| & \rightarrow 0 \\
\left\|T\left(x_{k}^{\prime}\right)-z\right\| & \rightarrow 0
\end{aligned}
$$

and hence

$$
\left\|T\left(x_{k}^{\prime}-x_{k}\right)-w\right\| \rightarrow 0
$$

as $k \rightarrow \infty$. Thus, since $\left\|x_{k}^{\prime}-x_{k}\right\|<2 \varepsilon$,

$$
\begin{equation*}
B\left(0, \delta_{0}\right) \subset \overline{T(B(0,2 \varepsilon))} \tag{8.4}
\end{equation*}
$$

Now we will show that the origin is interior to $T(B(0,2 \varepsilon))$. So fix $0<\varepsilon_{0}<\varepsilon$ and let $\left\{\varepsilon_{k}\right\}$ be a decreasing sequence of positive numbers such that $\sum_{k=1}^{\infty} \varepsilon_{k}<\varepsilon_{0}$. In view of (8.4), for each $k \geq 0$ there is a $\delta_{k}>0$ such that

$$
B\left(0, \delta_{k}\right) \subset \overline{T\left(B\left(0, \varepsilon_{k}\right)\right)}
$$

We may assume $\delta_{k} \rightarrow 0$ as $k \rightarrow \infty$. Fix $y \in B\left(0, \delta_{0}\right)$. Then $y$ is arbirarily close to elements of $T\left(B\left(0, \varepsilon_{0}\right)\right)$ and therefore there is an $x_{0} \in B\left(0, \varepsilon_{0}\right)$ such that

$$
\left\|y-T\left(x_{0}\right)\right\|<\delta_{1}
$$

i.e.,

$$
y-T\left(x_{0}\right) \in B\left(0, \delta_{1}\right) \subset \overline{T\left(B\left(0, \varepsilon_{1}\right)\right)}
$$

Thus there is an $x_{1} \in B\left(0, \varepsilon_{1}\right)$ such that

$$
\left\|y-T\left(x_{0}\right)-T\left(x_{1}\right)\right\|<\delta_{2}
$$

whence $y-T\left(x_{0}\right)-T\left(x_{1}\right) \in \overline{T\left(B\left(0, \varepsilon_{2}\right)\right)}$. By induction there is a sequence $\left\{x_{k}\right\}$ such that $x_{k} \in B\left(0, \varepsilon_{k}\right)$ and

$$
\left\|y-T\left(\sum_{k=0}^{m} x_{k}\right)\right\|<\delta_{m+1} .
$$

Thus $T\left(\sum_{k=0}^{m} x_{k}\right)$ converges to $y$ as $m \rightarrow \infty$. For all $m>0$, we have

$$
\sum_{k=0}^{m}\left\|x_{k}\right\|<\sum_{k=0}^{\infty} \varepsilon_{k}<\varepsilon_{0}
$$

which implies the series $\sum_{k=0}^{\infty} x_{k}$ converges absolutely; since $X$ is a Banach space and since $x_{k} \in B\left(0, \varepsilon_{0}\right)$ for all $k$, the series converges to an element $x$ in the closure of $B\left(0, \varepsilon_{0}\right)$ which is contained in $B\left(0,2 \varepsilon_{0}\right)$, see Proposition 8.10. The continuity of $T$ implies $T(x)=y$. Since $y$ is an arbitrary point in $B\left(0, \delta_{0}\right)$, we conclude that

$$
\begin{equation*}
B\left(0, \delta_{0}\right) \subset T\left(B\left(0,2 \varepsilon_{0}\right)\right) \subset T(B(0,2 \varepsilon)) \tag{8.5}
\end{equation*}
$$

which shows that the origin is interior to $T(B(0,2 \varepsilon)$.
Finally suppose $U$ is an open subset of $X$ and $y=T(x)$ for some $x \in U$. Let $\varepsilon>0$ be such that $B(x, \varepsilon) \subset U$. (8.5) states that there exists a $\delta>0$ such that

$$
B(0, \delta) \subset T(B(0, \varepsilon))
$$

From the linearity of $T$

$$
\begin{aligned}
T(B(x, \varepsilon)) & =\{T(x)+T(w) ; w \in B(0, \varepsilon)\} \\
& =\{y+T(w): w \in B(0, \varepsilon)\} \\
& \supset\{y+z: z \in B(0, \delta)\}=B(y, \delta)
\end{aligned}
$$

As an immediate consequence of the Open Mapping Theorem we have the following corollary.
8.25. Corollary. If $T$ is a one-to-one bounded linear mapping of a Banach space $X$ onto a Banach space $Y$, then $T^{-1}: Y \rightarrow X$ is a bounded linear mapping.

Proof. The existence and linearity of $T^{-1}$ are evident. If $U$ is an open subset of $X$, then the inverse image of $U$ under $T^{-1}$ is simply $T(U)$, which is open by Theorem 8.24. Thus, $T^{-1}$ is bounded, by Theorem 8.12.
8.26. Definition. The graph of a mapping $T: X \rightarrow Y$ is the set

$$
\{(x, T(x)): x \in X\} \subset X \times Y
$$

It is left as an exercise to show that the graph of a continuous linear mapping $T$ of a normed linear space $X$ into a normed linear space $Y$ is a closed subset of $X \times Y$. The following theorem shows that, when $X$ and $Y$ are Banach spaces, the converse is true, (cf. Exercise 5, Section 8.3).
8.27. Theorem (Closed Graph Theorem.). If $T$ is a linear mapping of a $B a$ nach space $X$ into a Banach space $Y$ and the graph of $T$ is closed in $X \times Y$, then $T$ is continuous.

Proof. For each $x \in X$ set

$$
\|x\|_{1}:=\|x\|+\|T(x)\| .
$$

It is readily verified that $\|\cdot\|_{1}$ is a norm on $X$. Let us show that it is complete. Suppose $\left\{x_{k}\right\}$ is a Cauchy sequence in $X$ with respect to $\|\cdot\|_{1}$. Then $\left\{x_{k}\right\}$ is a Cauchy sequence in $X$ and $\left\{T\left(x_{k}\right)\right\}$ is a Cauchy sequence in $Y$. Since $X$ and $Y$ are Banach spaces there exist $x \in X$ and $y \in Y$ such that $\left\|x_{k}-x\right\| \rightarrow 0$ and $\left\|T\left(x_{k}\right)-y\right\| \rightarrow 0$ as $k \rightarrow \infty$. Since the graph of $T$ is closed we must have

$$
(x, y)=(x, T(x))
$$

This implies $\left\|T\left(x_{k}\right)-T(x)\right\| \rightarrow 0$ as $k \rightarrow \infty$, and hence that $\left\|x-x_{k}\right\|_{1} \rightarrow 0$ as $k \rightarrow \infty$. Thus $X$ is a Banach space with respect to the norm $\|\cdot\|_{1}$.

Consider two copies of $X$, the first one with $X$ equipped with norm $\|\cdot\|_{1}$ and the second with $X$ equipped with $\|\cdot\|$. Then the identity mapping $I:\left(X,\|\cdot\|_{1}\right) \rightarrow$ $(X,\|\cdot\|)$ is continuous since $\|x\| \leq\|x\|_{1}$ for any $x \in X$. According to Corollary 8.25, $I$ is an open map which means that the inverse mapping of $I$ is also continuous. Hence, there is a constant $C$ such that

$$
\|x\|+\|T(x)\|=\|x\|_{1} \leq C\|x\|
$$

for all $x \in X$. Evidently $C \geq 1$. Thus

$$
\|T(x)\| \leq(C-1)\|x\|
$$

from which we conclude that $T$ is continuous.

## Exercises for Section 8.3

1. Suppose $T: X \rightarrow Y$ is a continuous linear mapping of a normed linear space $X$ into a normed linear space $Y$. Show that the graph of $T$ is closed in $X \times Y$.
2. Suppose $T: X \rightarrow Y$ is a univalent, continuous linear mapping of a Banach space $X$ into a Banach space $Y$. Prove that $T(X)$ is closed in $Y$ if and only if

$$
\|x\| \leq C\|T(x)\|
$$

for each $x \in X$.
3. Let $T_{k}$ be a sequence of bounded linear operators $T_{k}: X \rightarrow Y$ where $X$ is a Banach space and $Y$ is a normed linear space. If $\lim _{k \rightarrow \infty} T_{k}(x)$ exists for each $x \in X$, Corollary 8.22 yields the existence of a bounded linear operator $T: X \rightarrow Y$ such that $\lim _{k \rightarrow \infty} T_{k}(x)=T(x)$ for each $x \in X$. Give an example that shows that $T$ fails to be bounded if $X$ is not assumed to be a Banach space.
4. Let $M$ be an arbitrary closed subspace of a normed linear space $X$. Let us say that $x, y \in X$ are equivalent, written as $x \sim y$, if $x-y \in M$. We will denote by $[x]$ all elements $y \in X$ such that $x \sim y$.
(a) With $[x]+[y]:=[x+y]$ and $[c x]:=c[x]$ where $c \in \mathbb{R}$, prove that these operations are well-defined and that these cosets form a vector space.
(b) Let us define

$$
\|[x]\|:=\inf _{y \in M}\|x-y\|
$$

Prove that $\|[\cdot]\|$ is a norm on the space, $\mathcal{M}$, of all cosets $[x]$.
(c) The space $\mathcal{M}$ is called the quotient space and is denoted by $X / M$. Prove that if $M$ is a closed subspace of a Banach space $X$, then $X / M$ is also a Banach space.
5. Let $X$ denote the set of all sequences $\left\{a_{k}\right\}_{k=1}^{\infty}$ such that all but finitely many of the $a_{k}=0$.
(a) Show that $X$ is a linear space under the usual definitions of addition and scalar multiplication.
(b) Show that $\left\|\left\{a_{k}\right\}\right\|=\max _{k \geq 1}\left|a_{k}\right|$ is a norm on $X$.
(c) Define a mapping $T: X \rightarrow X$ by

$$
T\left(\left\{a_{k}\right\}\right)=\left\{k a_{k}\right\}
$$

Show that $T$ is a linear mapping.
(d) Show that the graph of $T$ is closed in $X \times X$.
(e) Show that $T$ is not continuous.

### 8.4. Dual Spaces

Here we introduce the important concept of the dual space of a normed linear space and the associated notion of weak topology.
8.28. Definition. Let $X$ be a normed linear space. The dual space $X^{*}$ of $X$ is the linear space of all bounded linear functionals on $X$ equipped with the norm

$$
\|f\|=\sup \{|f(x)|: x \in X,\|x\|=1\}
$$

In view of Theorem 8.16 we know that $X^{*}=\mathcal{B}(X, \mathbb{R})$ is a Banach space.

We will begin with a result concerning the relation between the topologies of $X$ and $X^{*}$. Recall that a topological space is separable if it contains a countable dense subset.
8.29. Theorem. $X$ is separable if $X^{*}$ is separable.

Proof. Let $\left\{f_{k}\right\}_{k=1}^{\infty}$ be a countable dense subset of $X^{*}$. For each $k$ there is an element $x_{k} \in X$ with $\left\|x_{k}\right\|=1$ such that

$$
\left|f_{k}\left(x_{k}\right)\right| \geq \frac{1}{2}\left\|f_{k}\right\|
$$

Let $W$ denote the set of all finite linear combinations of elements of $\left\{x_{k}\right\}$ with rational coefficients. Then it is easily verified that $\bar{W}$ is a subspace of $X$.

If $\bar{W} \neq X$, then there is an element $x_{0} \in X-\bar{W}$ and

$$
\inf _{w \in \bar{W}}\left\|x_{0}-w\right\|>0
$$

By Theorem 8.20 there exists an $f \in X^{*}$ such that

$$
f(w)=0 \text { for all } w \in \bar{W} \quad \text { and } \quad f\left(x_{0}\right)=1
$$

Since $\left\{f_{k}\right\}$ is dense in $X^{*}$ there is a subsequence $\left\{f_{k_{j}}\right\}$ for which

$$
\lim _{j \rightarrow \infty}\left\|f_{k_{j}}-f\right\|=0
$$

However, since $\left\|x_{k_{j}}\right\|=1$,

$$
\left\|f_{k_{j}}-f\right\| \geq\left\|f_{k_{j}}\left(x_{k_{j}}\right)-f\left(x_{k_{j}}\right)\right\|=\left\|f_{k_{j}}\left(x_{k_{j}}\right)\right\| \geq \frac{1}{2}\left\|f_{k_{j}}\right\|
$$

for each $j$. Thus $\left\|f_{k_{j}}\right\| \rightarrow 0$ as $j \rightarrow \infty$, which implies that $f=0$, contradicting the fact that $f\left(x_{0}\right)=1$. Thus $\bar{W}=X$.
8.30. Definition. If $X$ and $Y$ are linear spaces and $T: X \rightarrow Y$ is a one-to-one linear mapping of $X$ onto $Y$ we will call $T$ a linear isomorphism and say that $X$ and $Y$ are linearly isomorphic. If, in addition, $X$ and $Y$ are normed linear spaces and $\|T(x)\|=\|x\|$ for each $x \in X$, then $T$ is an isometric isomorphism and $X$ and $Y$ are isometrically isomorphic.

Denote by $X^{* *}$ the dual space of $X^{*}$. Suppose $X$ is normed linear space. For each $x \in X$ let $\Phi(x)$ be the linear functional on $X^{*}$ defined by

$$
\begin{equation*}
\Phi(x)(f)=f(x) \tag{8.6}
\end{equation*}
$$

for each $f \in X^{*}$. Since

$$
|\Phi(x)(f)| \leq\|f\|\|x\|
$$

the linear functional $\Phi(x)$ is bounded, in fact $\|\Phi(x)\| \leq\|x\|$. Thus $\Phi(x) \in X^{* *}$. It is readily verified that $\Phi$ is a bounded linear mapping of $X$ into $X^{* *}$ with $\|\Phi\| \leq 1$.

The following result is the key to understanding the relation between $X$ and $X^{*}$.
8.31. Proposition. Suppose $X$ is a normed linear space. Then

$$
\|x\|=\sup \left\{|f(x)|: f \in X^{*},\|f\|=1\right\}
$$

for each $x \in X$.
Proof. Fix $x \in X$. If $f \in X^{*}$ with $\|f\|=1$, then

$$
|f(x)| \leq\|f\|\|x\| \leq\|x\|
$$

If $x \neq 0$, then the distance from $x$ to the subspace $\{0\}$ is $\|x\|$ and according to Theorem 8.20 there is an element $g \in X^{*}$ such that $g(x)=1$ and $\|g\|=\frac{1}{\|x\|}$. Set $f=\|x\| g$, then $\|f\|=1$ and $f(x)=\|x\|$. Thus

$$
\sup \left\{|f(x)|: f \in X^{*},\|f\|=1\right\}=\|x\|
$$

for each $x \in X$.
8.32. Theorem. The mapping $\Phi$ is an isometric isomorphism of $X$ onto $\Phi(X)$.

Proof. In view of Proposition 8.31 we have

$$
\|\Phi(x)\|=\sup \left\{|f(x)|: f \in X^{*},\|f\|=1\right\}=\|x\|
$$

for each $x \in X$.
The mapping $\Phi$ is called the natural imbedding of $X$ in $X^{* *}$.
8.33. Definition. A normed linear space $X$ is said to be reflexive if

$$
\Phi(X)=X^{* *}
$$

in which case $X$ is isometrically isomorphic to $X^{* *}$.
Since $X^{* *}$ is a Banach space (see Theorem 8.16), it follows that every reflexive normed linear space is in fact a Banach space.
8.34. Examples. (i) The Banach space $\mathbb{R}^{n}$ is reflexive.
(ii) If $1 \leq p<\infty$ and $\frac{1}{p}+\frac{1}{p^{\prime}}=1$, then the linear mapping

$$
\Psi: L^{p^{\prime}}(X, \mu) \rightarrow\left(L^{p}(X, \mu)\right)^{*}
$$

defined by

$$
\Psi(g)(f)=\int_{X} g f d \mu
$$

for $g \in L^{p^{\prime}}(X, \mu)$ and $f \in L^{p}(X, \mu)$ is an isometric isomorphism of $L^{p^{\prime}}(X, \mu)$ onto $\left(L^{p}(X, \mu)\right)^{*}$. This is a rephrasing of Theorem 6.48 (note that for $p=1$, $\mu$ needs to be $\sigma$-finite.
(iii) For $1<p<\infty, L^{p}(X, \mu)$ is reflexive. In order to show this we fix $1<p<\infty$. We need to show that the natural imbedding $\Phi: L^{p}(X, \mu) \rightarrow\left(L^{p}(X, \mu)^{* *}\right.$ is on-to. From (ii) we have the isometric isomorphisms

$$
\begin{align*}
& \Psi_{1}: L^{p^{\prime}}(X, \mu) \rightarrow\left(L^{p}(X, \mu)\right)^{*}  \tag{8.7}\\
& \Psi_{1}(g)(f)=\int_{X} g f d \mu, f \in L^{p}
\end{align*}
$$

and

$$
\begin{gather*}
\Psi_{2}: L^{p}(X, \mu) \rightarrow\left(L^{p^{\prime}}(X, \mu)\right)^{*}  \tag{8.8}\\
\Psi_{2}(f)(g)=\int_{X} f g d \mu \quad, g \in L^{p^{\prime}}
\end{gather*}
$$

Let $w \in L^{p}(X, \mu)^{* *}$. From (8.7) we have $w \circ \Psi_{1} \in\left(L^{p^{\prime}}(X, \mu)\right)^{*}$. Thus, (8.8) implies that there exists $f \in L^{p}(X, \mu)$ such that $\Psi_{2}(f)=w \circ \Psi_{1}$. Therefore

$$
\begin{equation*}
\Psi_{2}(f)(g)=w \circ \Psi_{1}(g)=\int_{X} f g d \mu=\Psi_{1}(g)(f) \text { for all } g \in L^{p^{\prime}} \tag{8.9}
\end{equation*}
$$

We now proceed to check that

$$
\begin{equation*}
\Phi(f)=w \tag{8.10}
\end{equation*}
$$

Let $\alpha \in\left(L^{p}(X, \mu)\right)^{*}$. Then from (8.7) we obtain $g \in L^{p^{\prime}}(X, \mu)$ such that $\Psi_{1}(g)=\alpha$, and therefore from (8.9) we conclude

$$
\Phi(f)(\alpha)=\alpha(f)=\Psi_{1}(g)(f)=w\left(\Psi_{1}(g)\right)=w(\alpha)
$$

which proves (8.10).
(iv) Let $\Omega \subset \mathbb{R}^{n}$ be an open set. We recall that $\lambda$ denotes the Lebesgue measure in $\mathbb{R}^{n}$. We will prove that $L^{1}(\Omega, \lambda)$ is not reflexive. Proceeding by contradiction, if $L^{1}(\Omega, \lambda)$ were reflexive, and since $L^{1}(\Omega, \lambda)$ is separable (see Exercise 8 Section 8.4), it would follow that $L^{1}(\Omega, \lambda)^{* *}$ is separable, and hence $L^{1}(\Omega, \lambda)^{*}$ would also be separable. From (iii) we know that $L^{1}(\Omega, \lambda)^{*}$ is isometrically
isomorphic to $L^{\infty}(\Omega, \lambda)$. Therefore, we would conclude that $L^{\infty}(\Omega, \lambda)$ would be separable, which contradicts Exercise 9, Section 8.4.

In addition to the topology induced by the norm on a normed linear space it is useful to consider a smaller, i.e., "weaker" topology. The weak topology on a normed linear space $X$ is the smallest topology on $X$ with respect to which each $f \in X^{*}$ is continuous. That such a weak topology exists may be seen by observing that the intersection of any family of topologies for $X$ is a topology for $X$. In particular, the intersection of all topologies for $X$ that contains all sets of the form $f^{-1}(U)$ where $f \in X^{*}$ and $U$ is an open subset of $\mathbb{R}$ is precisely the weak topology. Any topology for $X$ with respect to which each $f \in X^{*}$ is continuous must contain the weak topology for $X$. Temporarily denote the weak topology by $\mathcal{T}_{w}$. Then $f^{-1}(U) \in \mathcal{T}_{w}$ whenever $f \in X^{*}$ and $U$ is open in $\mathbb{R}$. Consequently the family of all subsets of the form

$$
\left\{x:\left|f_{i}(x)-f_{i}\left(x_{0}\right)\right|<\varepsilon_{i}, 1 \leq i \leq m\right\}
$$

where $x_{0} \in X, m$ is a positive integer, and $\varepsilon_{i}>0, f_{i} \in X^{*}$ for $1 \leq i \leq m$ forms a base for $\mathcal{T}_{w}$. From this observation it is evident that a sequence $\left\{x_{k}\right\}$ in $X$ converges weakly (i.e., with respect to the topology $\mathcal{T}_{w}$ to $x \in X$ if and only if

$$
\lim _{k \rightarrow \infty} f\left(x_{k}\right)=f(x)
$$

for each $f \in X^{*}$.
In order to distinguish the weak topology from the topology induced by the norm, we will refer to the latter as the strong topology.
8.35. Theorem. Suppose $X$ is a normed linear space and the sequence $\left\{x_{k}\right\}$ converges weakly to $x \in X$. Then the following assertions hold:
(i) The sequence $\left\{\left\|x_{k}\right\|\right\}$ is bounded.
(ii) Let $W$ denote the subspace of $X$ spanned by $\left\{x_{k}: k=1,2, \ldots\right\}$. Then $x$ belongs to the closure of $W$, in the strong topology.
(iii)

$$
\|x\| \leq \liminf _{k \rightarrow \infty}\left\|x_{k}\right\|
$$

Proof. (i) Let $f \in X^{*}$. Since $\left\{f\left(x_{k}\right)\right\}$ is a convergent sequence in $\mathbb{R}$

$$
\sup \left\{\left|f\left(x_{k}\right)\right|: k=1,2, \ldots\right\}<\infty
$$

which may be written as

$$
\sup _{1 \leq k<\infty}\left|\Phi\left(x_{k}\right)(f)\right|<\infty
$$

Since this is true for each $f \in X^{*}$ and since $X^{*}$ is a Banach space (see Theorem 8.16), it follows from the Uniform Boundedness Principle, Theorem 8.21, that

$$
\sup _{1 \leq k<\infty}\left\|\Phi\left(x_{k}\right)\right\|<\infty
$$

In view of Theorem 8.32, this means

$$
\sup _{1 \leq k<\infty}\left\|x_{k}\right\|<\infty
$$

(ii) Let $\bar{W}$ denote the closure of $W$, in the strong topology. If $x \notin \bar{W}$, then by Theorem 8.20 there is an element $f \in X^{*}$ such that $f(x)=1$ and $f(w)=0$ for all $w \in \bar{W}$. But, as $f\left(x_{k}\right)=0$ for all $k$,

$$
f(x)=\lim _{k \rightarrow \infty} f\left(x_{k}\right)=0
$$

which contradicts the fact that $f(x)=1$. Thus $x \in \bar{W}$.
(iii) If $f \in X^{*}$ and $\|f\|=1$, then

$$
|f(x)|=\lim _{k \rightarrow \infty}\left|f\left(x_{k}\right)\right| \leq \liminf _{k \rightarrow \infty}\left\|x_{k}\right\|
$$

Since this is true for any such $f$,

$$
\|x\| \leq \liminf _{k \rightarrow \infty}\left\|x_{k}\right\|
$$

8.36. Theorem. If $X$ is a reflexive Banach space and $Y$ is a closed subspace of $X$, then $Y$ is a reflexive Banach space.

Proof. For any $f \in X^{*}$ let $f_{Y}$ denote the restriction of $f$ to $Y$. Then evidently $f_{Y} \in Y^{*}$ and $\left\|f_{Y}\right\| \leq\|f\|$. For $\omega \in Y^{* *}$ let $\omega_{Y}: X^{*} \rightarrow \mathbb{R}$ be given by

$$
\omega_{Y}(f)=\omega\left(f_{Y}\right)
$$

for each $f \in X^{*}$. Then $\omega_{Y}$ is a linear functional on $X^{*}$ and $\omega_{Y} \in X^{* *}$ since

$$
\left|\omega_{Y}(f)\right|=\left|\omega\left(f_{Y}\right)\right| \leq\|\omega\|\left\|f_{Y}\right\| \leq\|\omega\|\|f\|
$$

for each $f \in X^{*}$. Since $X$ is reflexive, there is an element $x_{0} \in X$ such that $\Phi\left(x_{0}\right)=\omega_{Y}$ where $\Phi$ is as in Definition 8.33 and therefore

$$
\omega_{Y}(f)=\Phi\left(x_{0}\right)(f)=f\left(x_{0}\right)
$$

for each $f \in X^{*}$. If $x_{0} \notin Y$, then by Theorem 8.20 there exists an $f \in X^{*}$ such that $f\left(x_{0}\right)=1$ and $f(y)=0$ for each $y \in Y$. This implies that $f_{Y}=0$ and hence we arrive at the contradiction

$$
1=f\left(x_{0}\right)=\omega_{Y}(f)=\omega\left(f_{Y}\right)=0
$$

Thus $x_{0} \in Y$.

For any $g \in Y^{*}$ there is, by Theorem 8.19, an $f \in X^{*}$ such that $f_{Y}=g$. Thus

$$
g\left(x_{0}\right)=f\left(x_{0}\right)=\omega_{Y}(f)=\omega\left(f_{Y}\right)=\omega(g)
$$

Thus the image of $x_{0}$ under the natural imbedding of $Y$ into $Y^{* *}$ is $\omega$. Since $\omega \in Y^{* *}$ is arbitrary, $Y$ is reflexive.
8.37. Theorem. If $X$ is a reflexive Banach space, then for any $M>0$, the closed ball $B:=\{x \in X:\|x\| \leq M\}$ is sequentially compact in the weak topology.

Proof. Assume first that $X$ is separable. Then since $X$ is reflexive, $X^{* *}$ is separable and, by Theorem $8.29, X^{*}$ is separable. Let $\left\{f_{m}\right\}_{m=1}^{\infty}$ be dense in $X^{*}$ and $\left\{x_{k}\right\} \in B$. Since $\left\{f_{1}\left(x_{k}\right)\right\}$ is bounded in $\mathbb{R}$ there be a subsequence $\left\{x_{k}^{1}\right\}$ of $\left\{x_{k}\right\}$ such that $\left\{f_{1}\left(x_{k}^{1}\right)\right\}$ converges in $\mathbb{R}$. Since the sequence $\left\{f_{2}\left(x_{k}^{1}\right)\right\}$ is bounded in $\mathbb{R}$ there is a subsequence $\left\{x_{k}^{2}\right\}$ of $\left\{x_{k}^{1}\right\}$ such that $\left\{f_{2}\left(x_{k}^{2}\right)\right\}$ converges in $\mathbb{R}$. Continuing in this way we obtain a sequence of subsequences of $\left\{x_{k}\right\}$ such that $\left\{x_{k}^{m}\right\}$ is a subsequence of $\left\{x_{k}^{m-1}\right\}$ for $m>1$ and $\left\{f_{m}\left(x_{k}^{m}\right)\right\}$ converges as $k \rightarrow \infty$ for each $m \geq 1$. Set $y_{k}=x_{k}^{k}$. Then $\left\{y_{k}\right\}$ is a subsequence of $\left\{x_{k}\right\}$ such that $\left\{f_{m}\left(y_{k}\right)\right\}$ converges as $k \rightarrow \infty$ for each $m \geq 1$. For arbitrary $f \in X^{*}$

$$
\begin{aligned}
\left|f\left(y_{k}\right)-f\left(y_{l}\right)\right| & \leq\left|f\left(y_{k}\right)-f_{m}\left(y_{k}\right)\right|+\left|f_{m}\left(y_{k}\right)-f_{m}\left(y_{l}\right)\right|+\left|f_{m}\left(y_{l}\right)-f\left(y_{l}\right)\right| \\
& \leq\left\|f_{m}-f\right\|\left(\left\|y_{k}\right\|+\left\|y_{l}\right\|\right)+\left|f_{m}\left(y_{k}\right)-f_{m}\left(y_{l}\right)\right|
\end{aligned}
$$

for any $k, l, m$. Given $\varepsilon>0$ there exists an $m$ such that

$$
\left\|f_{m}-f\right\|<\frac{\varepsilon}{4 M}
$$

where $M=\sup _{k \geq 1}\left\|x_{k}\right\|<\infty$. Since $\left\{f_{m}\left(y_{k}\right)\right\}$ is a Cauchy sequence, there is a positive integer $K$ such that

$$
\left|f_{m}\left(y_{k}\right)-f_{m}\left(y_{l}\right)\right|<\frac{\varepsilon}{2}
$$

whenever $k, l>K$. Thus, from the previous three inequalities,

$$
\left|f\left(y_{k}\right)-f\left(y_{l}\right)\right| \leq \frac{\varepsilon}{4 M} 2 M+\frac{\varepsilon}{2}<\varepsilon
$$

whenever $k, l>K$ and we see that $\left\{f\left(y_{k}\right)\right\}$ is a Cauchy sequence in $\mathbb{R}$. Set

$$
\alpha(f)=\lim _{k \rightarrow \infty} f\left(y_{k}\right)
$$

for each $f \in X^{*}$. Evidently $\alpha$ is a linear functional on $X^{*}$ and, since

$$
|\alpha(f)| \leq\|f\| M
$$

for each $f \in X^{*}$ we have $\alpha \in X^{* *}$. Since $X$ is reflexive, we know that the isometry $\Phi$ (see (8.6)) is onto $X^{* *}$. Thus, there exists $x \in X$ such that $\Phi(x)=\alpha$ and therefore

$$
\Phi(x)(f):=f(x)=\alpha(f)=\lim _{k \rightarrow \infty} f\left(y_{k}\right)
$$

for each $f \in X^{*}$. Thus $\left\{y_{k}\right\}$ converges to $x$ in the weak topology, and Theorem 8.35 gives $x \in B$.

Now suppose that $X$ is a reflexive Banach space, not necessarily separable, and suppose $\left\{x_{k}\right\} \in B$. It suffices to show that there exists $x \in B$ and a subsequence such that $\left\{x_{k}\right\} \rightarrow x$ weakly. Let $Y$ denote the closure in the strong topology of the subspace of $X$ spanned by $\left\{x_{k}\right\}$. Then $Y$ is obviously separable and, by Theorem 8.36, $Y$ is a reflexive Banach space. Thus there is a subsequence $\left\{x_{k_{j}}\right\}$ of $\left\{x_{k}\right\}$ that converges weakly in $Y$ to an element $x \in Y$, i.e.,

$$
g(x)=\lim _{j \rightarrow \infty} g\left(x_{k_{j}}\right)
$$

for each $g \in Y^{*}$. For any $f \in X^{*}$ let $f_{Y}$ denote the restriction of $f$ to $Y$. As in the proof of Theorem 8.36 we see that $f_{Y} \in Y^{*}$. Thus

$$
f(x)=f_{Y}(x)=\lim _{j \rightarrow \infty} f_{Y}\left(x_{k_{j}}\right)=\lim _{j \rightarrow \infty} f\left(x_{k_{j}}\right)
$$

Thus $\left\{x_{k_{j}}\right\}$ converges weakly to $x$ in $X$. Furthermore, since $\left\|x_{k}\right\| \leq 1$, the same is true for $x$ by Theorem 8.35 and so $x \in B$.
8.38. Example. Let $1<p<\infty$. Since $L^{p}\left(\mathbb{R}^{n}, \lambda\right)$ is a reflexive Banach space Theorem 8.37 implies that the ball

$$
B=\left\{f \in L^{p}:\|f\|_{p} \leq M\right\} \subset L^{p}\left(\mathbb{R}^{n}, \lambda\right)
$$

is sequentially compact in the weak topology. Thus, if $\left\{f_{k}\right\}$ is a sequence in $L^{p}\left(\mathbb{R}^{n}, \lambda\right)$ such that

$$
\left\|f_{k}\right\|_{p} \leq M \quad k=1,2,3 \ldots
$$

there exists a subsequence $\left\{f_{k_{j}}\right\}$ of $\left\{f_{k}\right\}$ and a function $f \in B$ such that

$$
\begin{equation*}
f_{k_{j}} \rightarrow f \text { weakly. } \tag{8.11}
\end{equation*}
$$

But (8.11) is equivalent to

$$
F\left(f_{k_{j}}\right) \rightarrow F(f) \text { for all } F \in L^{p}\left(\mathbb{R}^{n}, \lambda\right)^{*}
$$

Therefore, Examples 8.34 (ii) yields

$$
\int_{\mathbb{R}^{n}} f_{k_{j}} g d \lambda \rightarrow \int_{\mathbb{R}^{n}} f g d \lambda \quad \text { for all } g \in L^{p^{\prime}}\left(\mathbb{R}^{n}, \lambda\right)
$$

8.39. Definition. If $X$ is a normed linear space we may also consider the weak topology on $X^{*}$ i.e., the smallest topology on $X^{*}$ with respect to which each linear functional $\omega \in X^{* *}$ is continuous. It turns out to be convenient to consider an even weaker topology on $X^{*}$. The weak* topology on $X^{*}$ is defined as the smallest topology on $X^{*}$ with respect to which each linear functional $\omega \in \Phi(X) \subset X^{* *}$ is continuous. Here $\Phi$ is the natural imbedding of $X$ into $X^{* *}$. As in the case of the weak topology on $X$ we can, utilizing the natural imbedding, describe a base for the weak* topology on $X^{*}$ as the family of all sets of the form

$$
\left\{f \in X^{*}:\left|f\left(x_{i}\right)-f_{0}\left(x_{i}\right)\right|<\varepsilon_{i} \text { for } 1 \leq i \leq m\right\}
$$

where $m$ is any positive integer, $f_{0} \in X^{*}$, and $x_{i} \in X, \varepsilon_{i}>0$ for $1 \leq i \leq m$. Thus a sequence $\left\{f_{k}\right\}$ in $X^{*}$ converges in the weak* topology to an element $f \in X^{*}$ if and only if

$$
\lim _{k \rightarrow \infty} f_{k}(x)=f(x)
$$

for each $x \in X$. Of course, if $X$ is a reflexive Banach space, the weak and weak* topologies on $X^{*}$ coincide.

We remark that the base for the weak* topology on $X^{*}$ is similar in form to the base for the weak topology on $X$ except that the roles of $X$ and $X^{*}$ are interchanged.

The importance of the weak* topology is indicated by the following theorem.
8.40. Theorem (Alaoglu's Theorem). Suppose $X$ is a normed linear space. The unit ball $B:=\left\{f \in X^{*}:\|f\| \leq 1\right\}$ of $X^{*}$ is compact in the weak ${ }^{*}$ topology.

Proof. If $f \in B$, then $f(x) \in[-\|x\|,\|x\|]$ for each $x \in X$. Set $I_{x}=$ [ $-\|x\|,\|x\|]$ for $x \in X$. Then, according to Tychonoff's Theorem 3.43, the product

$$
P=\prod_{x \in X} I_{x}
$$

with the product topology, is compact. Recall that $B$ is by definition all functions $f$ defined on $X$ with the property that $f(x) \in I_{x}$ for each $x \in X$. Thus the set $B$ can be viewed as a subset $B^{\prime}$ of $P$. Moreover the relative topology induced on $B^{\prime}$ by the product topology is easily seen to coincide with the relative topology induced on $B$ by the weak* topology. Thus the proof will be complete if we show that $B^{\prime}$ is a closed subset of $P$ in the product topology.

Let $f$ be an element in the closure of $B^{\prime}$. Then given any $\varepsilon>0$ and $x \in X$ there is a $g \in B^{\prime}$ such that $|f(x)-g(x)|<\varepsilon$. Thus

$$
|f(x)| \leq \varepsilon+|g(x)| \leq \varepsilon+\|x\|
$$

since $g \in B^{\prime}$. Since $\varepsilon$ and $x$ are arbitrary,

$$
\begin{equation*}
|f(x)| \leq\|x\| \tag{8.12}
\end{equation*}
$$

for each $x \in X$.
Now suppose that $x, y \in X, \alpha, \beta \in \mathbb{R}$ and set $z=\alpha x+\beta y$. Then given any $\varepsilon>0$ there is a $g \in B^{\prime}$ such that

$$
|f(x)-g(x)|<\varepsilon, \quad|f(y)-g(y)|<\varepsilon, \quad|f(z)-g(z)|<\varepsilon
$$

Thus since $g$ is linear

$$
\begin{aligned}
\mid f(z) & -\alpha f(x)-\beta f(y) \mid \\
& \leq|f(z)-g(z)|+|\alpha||f(x)-g(x)|+|\beta||f(y)-g(y)| \\
& \leq \varepsilon(1+|\alpha|+|\beta|)
\end{aligned}
$$

from which it follows that

$$
f(\alpha x+\beta y)=\alpha f(x)+\beta f(y)
$$

i.e., $f$ is linear. In view of (8.12) $f \in B^{\prime}$. Thus $B^{\prime}$ is closed and hence compact in the product topology, from which it follows immediately that $B$ is compact in the weak* topology.

The proof of the following corollary of Theorem 8.40 is left as an exercise. Also, see Exercise 5, Section 8.4.
8.41. Corollary. The unit ball in a reflexive Banach space is both compact and sequentially compact in the weak topology.
8.42. Example. We apply Alaoglu's Theorem with $X=L^{1}\left(\mathbb{R}^{n}, \lambda\right)$, which is a normed linear space. Therefore, for any $M>0$, the ball $B=\left\{f \in L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}\right.$ : $\|f\| \leq M\} \subset L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}$ is compact in the weak* topology. We recall the isometric isomorphism

$$
\Psi: L^{\infty}\left(\mathbb{R}^{n}, \lambda\right) \rightarrow L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}
$$

given by

$$
\Psi(g)(f)=\int_{\mathbb{R}^{n}} g f, \quad f \in L^{1}\left(\mathbb{R}^{n}, \lambda\right)
$$

Using $\Psi$ we can rewrite the conclusion of Alaoglu's Theorem as

$$
B=\left\{\Psi(f) \in L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}:\|f\|_{\infty} \leq 1\right\} \subset L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}
$$

is compact in the weak* topology of $L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}$. From this it follows that if $\left\{f_{k}\right\} \subset$ $L^{\infty}\left(\mathbb{R}^{n}, \lambda\right)$ satisfies

$$
\left\|f_{k}\right\|_{\infty}=\left\|\Psi\left(f_{k}\right)\right\| \leq M, \quad k=1,2,3 \ldots
$$

then there exists a subsequence $f_{k_{j}}$ of $\left\{f_{k}\right\}$ and $f \in L^{\infty}\left(\mathbb{R}^{n}, \lambda\right)$ such that $\Psi\left(f_{k_{j}}\right) \rightarrow$ $\Psi(f)$ in the weak* topology of $L^{1}\left(\mathbb{R}^{n}, \lambda\right)^{*}$. This is equivalent to

$$
\Psi\left(f_{k_{j}}\right)(g) \rightarrow \Psi(f)(g) \quad \text { for all } g \in L^{1}\left(\mathbb{R}^{n}, \lambda\right)
$$

that is,

$$
\int_{\mathbb{R}^{n}} f_{k_{j}} g d \lambda \rightarrow \int_{\mathbb{R}^{n}} f g d \lambda \quad \text { for all } g \in L^{1}\left(\mathbb{R}^{n}, \lambda\right)
$$

## Exercises for Section 8.4

1. Suppose $X$ is a normed linear space and $\left\{f_{k}\right\}$ is a sequence in $X^{*}$ that converges in the weak* topology to $f \in X^{*}$. Show that

$$
\sup _{k \geq 1}\left\|f_{k}\right\|<\infty
$$

and

$$
\|f\| \leq \liminf _{k \rightarrow \infty}\left\|f_{k}\right\|
$$

2. Show that a subspace of a normed linear space is closed in the strong topology if and only if it is closed in the weak topology.
3. Prove that any subspace of a normed linear space cannot be open.
4. Show that a finite dimensional subspace of a normed linear space is closed in the strong topology.
5. Show that if $X$ is an infinite dimensional normed linear space, then there is a bounded sequence $\left\{x_{k}\right\}$ in $X$ of which no subsequence is convergent in the strong topology (Hint: Use Exercise 4, Section 8.4 and Exercise 1, Section 8.2). Thus conclude that the unit ball in any infinite dimensional normed linear space is not compact.
6. Show that any Banach space $X$ is isometrically isomorphic to a closed linear subspace of $C(\Gamma)$ (cf. Example $8.9(\mathrm{iv})$ ) where $\Gamma$ is a compact Hausdorff space. (Hint: Set

$$
\Gamma=\left\{f \in X^{*}:\|f\| \leq 1\right\}
$$

with the weak* topology. Use the natural imbedding of $X$ into $X^{* *}$.)
7. As usual, let $C[0,1]$ denote the space of continuous functions on $[0,1]$ endowed with the sup norm. Prove that if $f_{k}$ is a sequence of functions in $C[0,1]$ that converge weakly to $f$, then the sequence is bounded and $f_{k}(t) \rightarrow f(t)$ for each $t \in[0,1]$.
8. Suppose $\Omega$ is an open subset of $\mathbb{R}^{n}$, and let $\lambda$ denote Lebesgue measure on $\Omega$. Set

$$
\mathcal{P}=\left\{x=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathbb{R}^{n}: x_{j} \quad \text { is rational for each } \quad 1 \leq j \leq n\right\}
$$

and let $\mathcal{Q}$ denote the set of all open cubes in $\mathbb{R}^{n}$ with edges parallel to the coordinate axes and vertices in $\mathcal{P}$. (i) Show that if $E$ is a Lebesgue measurable subset of $\mathbb{R}^{n}$ with $\lambda(E)<\infty$ and $\varepsilon>0$, then there exists a disjoint finite sequence $\left\{Q_{k}\right\}_{k=1}^{m}$ with each $Q_{k} \in \mathcal{Q}$ such that

$$
\left\|\chi_{E}-\sum_{k=1}^{m} \chi_{Q_{k}}\right\|_{L^{p}(\Omega, \lambda)}<\varepsilon
$$

(ii) Show that the set of all finite linear combinations of elements of $\left\{\chi_{Q}: Q \in \mathcal{Q}\right\}$ with rational coefficients is dense in $L^{p}(\Omega, \lambda)$.
(iii) Conclude that $L^{p}(\Omega, \lambda)$ is separable.
9. Suppose $\Omega$ is an open subset of $\mathbb{R}^{n}$ and let $\lambda$ denote Lebesgue measure on $\Omega$. Show that $L^{\infty}(\Omega, \lambda)$ is not separable. (Hint: If $B(x, r) \in \Omega$ and $0<r_{1}<r_{2} \leq r$, then

$$
\left.\left\|\chi_{B\left(x, r_{1}\right)}-\chi_{B\left(x, r_{2}\right)}\right\|_{L^{\infty}(\Omega, \lambda)}=1 .\right)
$$

10. Referring to Exercise 2, Section 8.1, prove that there is a natural isomorphism between $X^{*}$ and $\prod_{i=1}^{m} X_{i}^{*}$. Thus conclude that $X$ is reflexive if each $X_{i}$ is reflexive.
11. Let $X$ be a normed linear space and suppose that the sequence $\left\{x_{k}\right\}$ converges to $x \in X$ in the strong topology. Show that $\left\{x_{k}\right\}$ converges weakly to $x$.

### 8.5. Hilbert Spaces

We consider in this section Hilbert spaces, i.e., Banach spaces in which the norm is induced by an inner product. This additional structure allows us to study the representation of elements of the space in terms of orthonormal systems.
8.43. Definition. An inner product on a linear space $X$ is a real-valued function $(x, y) \mapsto\langle x, y\rangle$ on $X \times X$ such that for $x, y, z \in X$ and $\alpha, \beta \in \mathbb{R}$

$$
\begin{aligned}
\langle x, y\rangle & =\langle y, x\rangle \\
\langle\alpha x+\beta y, z\rangle & =\alpha\langle x, z\rangle+\beta\langle y, z\rangle \\
\langle x, x\rangle & \geq 0 \\
\langle x, x\rangle & =0 \text { if, and only if, } x=0
\end{aligned}
$$

8.44. Theorem. Suppose that $X$ is a linear space on which an inner product $\langle\cdot, \cdot\rangle$ is defined. We can define a norm on $X$ as follows

$$
\|x\|=\sqrt{\langle x, x\rangle} .
$$

Proof. It follows immediately from the definition of an inner product that

$$
\begin{aligned}
\|x\| & \geq 0 \text { for all } x \in X \\
\|x\| & =0 \text { if, and only if } x=0 \\
\|\alpha x\| & =|\alpha|\|x\| \text { for all } \alpha \in \mathbb{R}, x \in X
\end{aligned}
$$

Only the triangle inequality remains to be proved. To do this, we first prove the Schwarz Inequality

$$
|\langle x, y\rangle| \leq\|x\|\|y\|
$$

Suppose $x, y \in X$ and $\lambda \in \mathbb{R}$. From the properties of an inner product,

$$
\begin{aligned}
0 \leq\|x-\lambda y\|^{2} & =\langle x-\lambda y, x-\lambda y\rangle \\
& =\|x\|^{2}-2 \lambda\langle x, y\rangle+\lambda^{2}\|y\|^{2}
\end{aligned}
$$

and thus

$$
2\langle x, y\rangle \leq \frac{1}{\lambda}\|x\|^{2}+\lambda\|y\|^{2}
$$

for any $\lambda>0$. Assuming $y \neq 0$ and setting $\lambda=\frac{\|x\|}{\|y\|}$, we see that

$$
\begin{equation*}
\langle x, y\rangle \leq\|x\|\|y\| . \tag{8.13}
\end{equation*}
$$

Note that (8.13) also holds in case $y=0$. Since

$$
-\langle x, y\rangle=\langle x,-y\rangle \leq\|x\|\|y\|
$$

we see that

$$
\begin{equation*}
|\langle x, y\rangle| \leq\|x\|\|y\| \tag{8.14}
\end{equation*}
$$

for all $x, y \in X$.
For the triangle inequality, observe that

$$
\begin{aligned}
\|x+y\|^{2} & =\|x\|^{2}+2\langle x, y\rangle+\|y\|^{2} \\
& \leq\|x\|^{2}+2\|x\|\|y\|+\|y\|^{2} \\
& =(\|x\|+\|y\|)^{2} .
\end{aligned}
$$

Thus

$$
\|x+y\| \leq\|x\|+\|y\|
$$

for any $x, y \in X$, from which we see that $\|\cdot\|$ is a norm on $X$.

Thus we see that a linear space equipped with an inner product is a normed linear space. The inequality (8.14) is called the Schwarz Inequality.
8.45. Definition. A Hilbert space is a linear space with an inner product that is a Banach space with respect to the norm induced by the inner product (as in Theorem 8.44).
8.46. Definition. We will say that two elements $x, y$ in a Hilbert space $H$ are orthogonal if $\langle x, y\rangle=0$. If $M$ is a subspace of $H$, we set

$$
M^{\perp}=\{x \in H:\langle x, y\rangle=0 \text { for all } y \in M\}
$$

It is easily seen that $M^{\perp}$ is a subspace of $H$.
We next investigate the "geometry" of a Hilbert space.
8.47. Theorem. Suppose $M$ is a closed subspace of a Hilbert space $H$. Then for each $x_{0} \in H$ there exists a unique $y_{0} \in M$ such that

$$
\left\|x_{0}-y_{0}\right\|=\inf _{y \in M}\left\|x_{0}-y\right\|
$$

Moreover $y_{0}$ is the unique element of $M$ such that $x_{0}-y_{0} \in M^{\perp}$.
Proof. If $x_{0} \in M$ the assertion is obvious, so assume $x_{0} \in H-M$. Since $M$ is closed and $x_{0} \notin M$,

$$
d=\inf _{y \in M}\left\|x_{0}-y\right\|>0
$$

There is a sequence $\left\{y_{k}\right\}$ in $M$ such that

$$
\lim _{k \rightarrow \infty}\left\|x_{0}-y_{k}\right\|=d
$$

For any $k, l$

$$
\begin{gathered}
\left\|\left(x_{0}-y_{k}\right)-\left(x_{0}-y_{l}\right)\right\|^{2}+\left\|\left(x_{0}-y_{k}\right)+\left(x_{0}-y_{l}\right)\right\|^{2} \\
=2\left\|x_{0}-y_{k}\right\|^{2}+2\left\|x_{0}-y_{l}\right\|^{2} \\
\left\|y_{k}-y_{l}\right\|^{2}=2\left(\left\|x_{0}-y_{k}\right\|^{2}+\left\|x_{0}-y_{l}\right\|^{2}\right)-4\left\|x_{0}-\frac{1}{2}\left(y_{k}+y_{l}\right)\right\|^{2} . \\
\leq 2\left(\left\|x_{0}-y_{k}\right\|^{2}+\left\|x_{0}-y_{l}\right\|^{2}\right)-4 d^{2} .
\end{gathered}
$$

Since the right hand side of the last inequality above tends to 0 as $k, l \rightarrow \infty$ we see that $\left\{y_{k}\right\}$ is a Cauchy sequence in $H$ and consequently converges to an element $y_{0}$. Since $M$ is closed, $y_{0} \in M$.

Now let $y \in M, \lambda \in \mathbb{R}$ and compute

$$
\begin{aligned}
d^{2} & \leq\left\|x_{0}-\left(y_{0}+\lambda y\right)\right\|^{2} \\
& =\left\|x_{0}-y_{0}\right\|^{2}-2 \lambda\left\langle x_{0}-y_{0}, y\right\rangle+\lambda^{2}\|y\|^{2} \\
& =d^{2}-2 \lambda\left\langle x_{0}-y_{0}, y\right\rangle+\lambda^{2}\|y\|^{2}
\end{aligned}
$$

whence

$$
\left\langle x_{0}-y_{0}, y\right\rangle \leq \frac{\lambda}{2}\|y\|^{2}
$$

for any $\lambda>0$. Since $\lambda$ is otherwise arbitrary we conclude that

$$
\left\langle x_{0}-y_{0}, y\right\rangle \leq 0
$$

for each $y \in M$. But then

$$
-\left\langle x_{0}-y_{0}, y\right\rangle=\left\langle x_{0}-y_{0},-y\right\rangle \leq 0
$$

for each $y \in M$ and thus $\left\langle x_{0}-y_{0}, y\right\rangle=0$ for each $y \in M$, i.e., $x_{0}-y_{0} \in M^{\perp}$.
If $y_{1} \in M$ is such that

$$
\left\|x_{0}-y_{1}\right\|=\inf _{y \in M}\left\|x_{0}-y\right\|
$$

then the above argument shows that $x_{0}-y_{1} \in M^{\perp}$. Hence, $y_{1}-y_{0}=\left(x_{0}-y_{0}\right)-$ $\left(x_{0}-y_{1}\right) \in M^{\perp}$. Since we also have $y_{1}-y_{0} \in M$, this implies

$$
\left\|y_{1}-y_{0}\right\|^{2}=\left\langle y_{1}-y_{0}, y_{1}-y_{0}\right\rangle=0
$$

i.e., $y_{1}=y_{0}$.
8.48. Theorem. Suppose $M$ is a closed subspace of a Hilbert space $H$. Then for each $x \in H$ there exists a unique pair of elements $y \in M$ and $z \in M^{\perp}$ such that $x=y+z$.

Proof. We may assume that $M \neq H$. Let $x \in H$. According to Theorem 8.47 there is a $y \in M$ such that $z=x-y \in M^{\perp}$. This establishes the existence of $y \in M, z \in M^{\perp}$ such that $x=y+z$.

To show uniqueness, suppose that $y_{1}, y_{2} \in M$, and $z_{1}, z_{2} \in M^{\perp}$ are such that $y_{1}+z_{1}=y_{2}+z_{2}$. Then

$$
y_{1}-y_{2}=z_{2}-z_{1}
$$

which means that $y_{1}-y_{2} \in M \cap M^{\perp}$. Thus

$$
\left\|y_{1}-y_{2}\right\|^{2}=\left\langle y_{1}-y_{2}, y_{1}-y_{2}\right\rangle=0
$$

whence $y_{1}=y_{2}$, This in turn implies that $z_{1}=z_{2}$.

If $y \in H$ is fixed and we define the function

$$
f(x)=\langle y, x\rangle
$$

for $x \in H$, then $f$ is a linear functional on $H$. Furthermore from Schwarz's inequality

$$
|f(x)|=|\langle y, x\rangle| \leq\|y\|\|x\|
$$

which implies that $\|f\| \leq\|y\|$. If $y=0$, then $\|f\|=0$. If $y \neq 0$, then

$$
f\left(\frac{y}{\|y\|}\right)=\|y\|
$$

Thus $\|f\|=\|y\|$. Using Theorem 8.48, we will show that every continuous linear functional on $H$ is of this form.
8.49. Theorem (Riesz Representation Theorem). Suppose $H$ is a Hilbert space. Then for each $f \in H^{*}$ there exists a unique $y \in H$ such that

$$
\begin{equation*}
f(x)=\langle y, x\rangle \tag{8.15}
\end{equation*}
$$

for each $x \in H$. Moreover, under this correspondence $H$ and $H^{*}$ are isometrically isomorphic.

Proof. Suppose $f \in H^{*}$. If $f=0$, then (8.15) holds with $y=0$. So assume that $f \neq 0$. Then

$$
M=\{x \in H: f(x)=0\}
$$

is a closed subspace of $H$ and $M \neq H$. We infer from Theorem 8.48 that there is an element $x_{0} \in M^{\perp}$ with $x_{0} \neq 0$. Since $x_{0} \notin M, f\left(x_{0}\right) \neq 0$. Since for each $x \in H$

$$
f\left(x-\frac{f(x)}{f\left(x_{0}\right)} x_{0}\right)=0
$$

we see that

$$
x-\frac{f(x)}{f\left(x_{0}\right)} x_{0} \in M
$$

for each $x \in H$. Thus

$$
\left\langle x-\frac{f(x)}{f\left(x_{0}\right)} x_{0}, x_{0}\right\rangle=0
$$

for each $x \in H$. This last equation may be rewritten as

$$
f(x)=\frac{f\left(x_{0}\right)}{\left\|x_{0}\right\|^{2}}\left\langle x_{0}, x\right\rangle
$$

Thus if we set $y=\frac{f\left(x_{0}\right)}{\left\|x_{0}\right\|^{2}} x_{0}$ we see that (8.15) holds. We have already observed that the norm of a linear functional $f$ satisfying (8.15) is $\|y\|$.

If $y_{1}, y_{2} \in H$ are such that $\left\langle y_{1}, x\right\rangle=\left\langle y_{2}, x\right\rangle$ for all $x \in H$, then

$$
\left\langle y_{1}-y_{2}, x\right\rangle=0
$$

for all $x \in X$. Thus, in particular,

$$
\left\|y_{1}-y_{2}\right\|^{2}=\left\langle y_{1}-y_{2}, y_{1}-y_{2}\right\rangle=0
$$

whence $y_{1}=y_{2}$. This shows that the $y$ that represents $f$ in (8.15) is unique.
We may rephrase the above results as follows. Let $\Psi: H \rightarrow H^{*}$ be defined for each $x \in H$ by

$$
\Psi(x)(y)=\langle x, y\rangle
$$

for all $y \in H$. Then $\Psi$ is a one-to-one linear mapping of $H$ onto $H^{*}$. Furthermore $\|\Psi(x)\|=\|x\|$ for each $x \in H$. Thus $\Psi$ is an isometric isomorphism of $H$ onto $H^{*}$.
8.50. Theorem. Any Hilbert space $H$ is a reflexive Banach space and consequently the set $\{x \in H:\|x\| \leq 1\}$ is compact in the weak topology.

Proof. We first show that $H^{*}$ is a Hilbert space. Let $\Psi$ be as in the proof of Theorem 8.49 above and define

$$
\begin{equation*}
\langle f, g\rangle=\left\langle\Psi^{-1}(f), \Psi^{-1}(g)\right\rangle \tag{8.16}
\end{equation*}
$$

for each pair $f, g \in H^{*}$. Note that the right member of the equation above is the inner product in $H$. That (8.16) defines an inner product on $H^{*}$ follows immediately from the properties of $\Psi$. Furthermore

$$
\langle f, f\rangle=\left\langle\Psi^{-1}(f), \Psi^{-1}(f)\right\rangle=\|f\|^{2}
$$

for each $f \in H^{*}$. Thus the norm on $H^{*}$ is induced by this inner product. If $\omega \in H^{* *}$, then by Theorem 8.49, there is an element $g \in H^{*}$ such that

$$
\omega(f)=\langle g, f\rangle
$$

for each $f \in H^{*}$. Again by Theorem 8.49 there is an element $x \in H$ such that $\Psi(x)=g$. Thus

$$
\omega(f)=\langle g, f\rangle=\langle\Psi(x), f\rangle=\left\langle x, \Psi^{-1}(f)\right\rangle=f(x)
$$

for each $f \in H^{*}$ from which we conclude that $H$ is reflexive. Thus $\{x \in H:\|x\| \leq$ $1\}$ is compact in the weak topology by Theorem 8.41.

We next consider the representation of elements of a Hilbert space by "Fourier series."
8.51. Definition. A subset $\mathcal{F}$ of a Hilbert space is an orthonormal family if for each pair of elements $x, y \in \mathcal{F}$

$$
\begin{aligned}
& \langle x, y\rangle=0 \text { if } x \neq y \\
& \langle x, y\rangle=1 \text { if } x=y .
\end{aligned}
$$

An orthonormal family $\mathcal{F}$ in $H$ is complete if the only element $x \in H$ for which $\langle x, y\rangle=0$ for all $y \in \mathcal{F}$ is $x=0$.
8.52. Theorem. Every Hilbert space contains a complete orthonormal family.

Proof. This assertion follows from the Hausdorff Maximal Principle (see Exercise 1, Section 8.5).
8.53. Theorem. If $H$ is a separable Hilbert space and $\mathcal{F}$ is an orthonormal family in $H$, then $\mathcal{F}$ is at most countable.

Proof. If $x, y \in \mathcal{F}$ and $x \neq y$, then

$$
\|x-y\|^{2}=\|x\|^{2}-2\langle x, y\rangle+\|y\|^{2}=2
$$

Thus

$$
B\left(x, \frac{1}{2}\right) \bigcap B\left(y, \frac{1}{2}\right)=\emptyset
$$

whenever $x, y$ are distinct elements of $\mathcal{F}$. If $\mathcal{E}$ is a countable dense subset of $H$, then for each $x \in \mathcal{F}$ the set $B\left(x, \frac{1}{2}\right)$ must contain an element of $\mathcal{E}$.

We next study the properties of orthonormal families beginning with countable orthonormal families.
8.54. Theorem. Suppose $\left\{x_{k}\right\}_{k=1}^{\infty}$ is an orthonormal sequence in a Hilbert space $H$. Then the following assertions hold:
(i) For each $x \in H$

$$
\sum_{k=1}^{\infty}\left\langle x, x_{k}\right\rangle^{2} \leq\|x\|^{2}
$$

(ii) If $\left\{\alpha_{k}\right\}$ is a sequence of real numbers, then

$$
\left\|x-\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k}\right\| \leq\left\|x-\sum_{k=1}^{m} \alpha_{k} x_{k}\right\|
$$

for each $m \geq 1$.
(iii) If $\left\{\alpha_{k}\right\}$ is a sequence of real numbers then $\sum_{k=1}^{\infty} \alpha_{k} x_{k}$ converges in $H$ if, and only if, $\sum_{k=1}^{\infty} \alpha_{k}^{2}$ converges in $\mathbb{R}$, in which case the sum is independent of the order in which the terms are arranged, i.e., the series converges unconditionally, and

$$
\left\|\sum_{k=1}^{\infty} \alpha_{k} x_{k}\right\|^{2}=\sum_{k=1}^{\infty} \alpha_{k}^{2}
$$

Proof. (i) For any positive integer $m$

$$
\begin{aligned}
0 & \leq\left\|x-\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k}\right\|^{2} \\
& =\|x\|^{2}-2\left\langle x, \sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k}\right\rangle+\left\|\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k}\right\|^{2} \\
& =\|x\|^{2}-2 \sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle^{2}+\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle^{2} \\
& =\|x\|^{2}-\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle^{2} .
\end{aligned}
$$

Thus

$$
\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle^{2} \leq\|x\|^{2}
$$

for any $m \geq 1$ from which assertion (i) follows.
(ii) Fix a positive integer $m$ and let $M$ denote the subspace of $H$ spanned by $\left\{x_{1}, x_{2}, \ldots, x_{m}\right\}$. Then $M$ is finite dimensional and hence closed; see Exercise 4, Section 8.4. Since for any $1 \leq k \leq m$

$$
\left\langle x_{k}, x-\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k}\right\rangle=0
$$

we see that

$$
x-\sum_{k=1}^{m}\left\langle x, x_{k}\right\rangle x_{k} \in M^{\perp}
$$

In view of Theorem 8.47 this implies assertion (ii) since

$$
\sum_{k=1}^{m} \alpha_{k} x_{k} \in M
$$

(iii) For any positive integers $m$ and $l$ with $m>l$

$$
\left\|\sum_{k=1}^{m} \alpha_{k} x_{k}-\sum_{k=1}^{l} \alpha_{k} x_{k}\right\|^{2}=\left\|\sum_{k=l+1}^{m} \alpha_{k} x_{k}\right\|^{2}=\sum_{k=l+1}^{m} \alpha_{k}^{2}
$$

Thus the sequence $\left\{\sum_{k=1}^{m} \alpha_{k} x_{k}\right\}_{m=1}^{\infty}$ is a Cauchy sequence in $H$ if and only if $\sum_{k=1}^{\infty} \alpha_{k}^{2}$ converges in $\mathbb{R}$.

Suppose $\sum_{k=1}^{\infty} \alpha_{k}^{2}<\infty$. Since for any $m$

$$
\left\|\sum_{k=1}^{m} \alpha_{k} x_{k}\right\|^{2}=\sum_{k=1}^{m} \alpha_{k}^{2}
$$

we see that

$$
\left\|\sum_{k=1}^{\infty} \alpha_{k} x_{k}\right\|^{2}=\sum_{k=1}^{\infty} \alpha_{k}^{2}
$$

Let $\left\{\alpha_{k_{j}}\right\}$ be any rearrangement of the sequence $\left\{\alpha_{k}\right\}$. Then for any $m$

$$
\begin{equation*}
\left\|\sum_{k=1}^{m} \alpha_{k} x_{k}-\sum_{j=1}^{m} \alpha_{k_{j}} x_{k_{j}}\right\|^{2}=\sum_{k=1}^{m} \alpha_{k}^{2}-2 \sum_{k_{j} \leq m} \alpha_{k_{j}}^{2}+\sum_{j=1}^{m} \alpha_{k_{j}}^{2} \tag{8.17}
\end{equation*}
$$

Since the sum of the series $\sum_{k=1}^{\infty} \alpha_{k}^{2}$ is independent of the order of the terms, the last member of (8.17) converges to 0 as $m \rightarrow \infty$.
8.55. Theorem. Suppose $\mathcal{F}$ is an orthonormal system in a Hilbert space $H$ and $x \in H$. Then
(i) The set $\{y \in \mathcal{F}:\langle x, y\rangle \neq 0\}$ is at most countable,
(ii) The series $\sum_{y \in \mathcal{F}}\langle x, y\rangle y$ converges unconditionally in $H$.

Proof. (i) Let $\varepsilon>0$ and set

$$
\mathcal{F}_{\varepsilon}=\{y \in \mathcal{F}:|\langle x, y\rangle|>\varepsilon\}
$$

In view of Theorem $8.54(\mathrm{a})$, the number of elements in $\mathcal{F}_{\varepsilon}$ cannot exceed $\frac{|x|^{2}}{\varepsilon^{2}}$ and thus $\mathcal{F}_{\varepsilon}$ is a finite set. Since

$$
\{y \in \mathcal{F}:\langle x, y\rangle \neq 0\}=\bigcup_{k=1}^{\infty}\left\{y \in \mathcal{F}:|\langle x, y\rangle|>\frac{1}{k}\right\}
$$

we see that $\{y \in \mathcal{F}:\langle x, y\rangle \neq 0\}$ is at most countable. (ii) Set $\left\{y_{k}\right\}_{k=1}^{\infty}=\{y \in$
$\mathcal{F}:\langle x, y\rangle \neq 0\}$. Then, according to Theorem 8.54(i), (iii), the series $\sum_{k=1}^{\infty}\left\langle x, y_{k}\right\rangle y_{k}$ converges unconditionally in $H$.
8.56. Theorem. Suppose $\mathcal{F}$ is a complete orthonormal system in a Hilbert space $H$. Then for each $x \in H$

$$
x=\sum_{y \in \mathcal{F}}\langle x, y\rangle y
$$

where all but countably many terms in the series are equal to 0 and the series converges unconditionally in $H$.

Proof. Let $Y$ denote the subspace of $H$ spanned by $\mathcal{F}$ and let $M$ denote the closure of $Y$. If $z \in M^{\perp}$, then $\langle z, y\rangle=0$ for each $y \in \mathcal{F}$. Since $\mathcal{F}$ is complete this implies that $z=0$. Thus $M^{\perp}=\{0\}$.

Fix $x \in H$. By Theorem 8.55 (ii), the series $\sum_{y \in \mathcal{F}}\langle x, y\rangle y$ converges unconditionally to an element $\bar{x} \in H$.

Set $\left\{y_{k}\right\}_{k=1}^{\infty}=\{y \in \mathcal{F}:\langle x, y\rangle \neq 0\}$. If $y \in \mathcal{F}$ and $\langle x, y\rangle=0$, then

$$
\left\langle x-\sum_{k=1}^{m}\left\langle x, y_{k}\right\rangle y_{k}, y\right\rangle=0
$$

for each $m \geq 1$. Then

$$
\begin{aligned}
|\langle x-\bar{x}, y\rangle| & =\left|\left\langle x-\sum_{k=1}^{m}\left\langle x, y_{k}\right\rangle y_{k}, y\right\rangle-\langle x-\bar{x}, y\rangle\right| \\
& =\left|\left\langle\bar{x}-\sum_{k=1}^{m}\left\langle x, y_{k}\right\rangle y_{k}, y\right\rangle\right| \\
& \leq\left\|\bar{x}-\sum_{k=1}^{m}\left\langle x, y_{k}\right\rangle y_{k}\right\|\|y\|
\end{aligned}
$$

and we see that $\langle x-\bar{x}, y\rangle=0$.
If $1 \leq l \leq m$, then

$$
\left\langle x-\sum_{k=1}^{m}\left\langle x, y_{k}\right\rangle y_{k}, y_{l}\right\rangle=0
$$

and hence $\left\langle x-\bar{x}, y_{l}\right\rangle=0$ for every $l \geq 1$.
We have thus shown that

$$
\begin{equation*}
\langle x-\bar{x}, y\rangle=0 \tag{8.18}
\end{equation*}
$$

for each $y \in \mathcal{F}$ from which it follows immediately that (8.18) holds for each $y \in Y$.
If $w \in M$, then there is a sequence $\left\{w_{k}\right\}$ in $Y$ such that $\left\|w-w_{k}\right\| \rightarrow 0$ as $k \rightarrow \infty$. Thus

$$
\langle x-\bar{x}, w\rangle=\lim _{k \rightarrow \infty}\left\langle x-\bar{x}, w_{k}\right\rangle=0
$$

which means that $x-\bar{x} \in M^{\perp}=\{0\}$.
8.57. Examples. (i) The Banach space $\mathbb{R}^{n}$ is a Hilbert space with inner product

$$
\left\langle\left(x_{1}, x_{2}, \ldots, x_{n}\right),\left(y_{1}, y_{2}, \ldots, y_{n}\right)\right\rangle=x_{1} y_{1}+x_{2} y_{2}+\cdots+x_{n} y_{n}
$$

(ii) The Banach space $L^{2}(X, \mu)$ is a Hilbert space with inner product

$$
\langle f, g\rangle=\int_{X} f g d \mu
$$

(iii) The Banach space $l^{2}$ is a Hilbert space with inner product

$$
\left\langle\left\{a_{k}\right\},\left\{b_{k}\right\}\right\rangle=\sum_{k=1}^{\infty} a_{k} b_{k}
$$

Suppose $H$ is a separable Hilbert space containing a countable orthonormal system $\left\{x_{k}\right\}_{k=1}^{\infty}$. For any $x \in H$ the sequence $\left\{\left\langle x, x_{k}\right\rangle\right\} \in l^{2}$ and

$$
\sum_{k=1}^{\infty}\left\langle x, x_{k}\right\rangle^{2}=\|x\|^{2}
$$

On the other hand, if $\left\{a_{k}\right\} \in l^{2}$, then according to Theorem 8.54 the series $\sum_{k=1}^{\infty} a_{k} x_{k}$ converges in $H$. Set $x=\sum_{k=1}^{\infty} a_{k} x_{k}$. Then

$$
\left\langle x, x_{l}\right\rangle=\lim _{m \rightarrow \infty}\left\langle\sum_{k=1}^{m} a_{k} x_{k}, x_{l}\right\rangle=a_{l}
$$

for each $l$ and hence

$$
\sum_{k=1}^{\infty} a_{k}^{2}=\|x\|^{2}
$$

Thus the linear mapping $T: H \rightarrow l^{2}$ given by

$$
T(x)=\left\{\left\langle x, x_{k}\right\rangle\right\}
$$

is an isometric isomorphism of $H$ onto $l^{2}$.
If $\Omega$ is a open subset of $\mathbb{R}^{n}$ and $\lambda$ denotes Lebesgue measure on $\Omega$ then $L^{2}(\Omega, \lambda)$ is separable and hence isometrically isomorphic to $l^{2}$.

## Exercises for Section 8.5

1. Show that every Hilbert space contains a complete orthonormal system. (Hint:

Observe that an orthonormal system in a Hilbert space $H$ is complete if and only if it is maximal with respect to set inclusion i.e., it is not contained in any other orthonormal system.)
2. Suppose $T$ is a linear mapping of a Hilbert space $H$ into a Hilbert space $E$ such that $\|T(x)\|=\|x\|$ for each $x \in H$. Show that

$$
\langle T(x), T(y)\rangle=\langle x, y\rangle
$$

for any $x, y \in H$.
3. Show that a Hilbert space that contains a countable complete orthonormal system is separable.
4. Fix $1<p<\infty$. Set $x^{m}=\left\{x_{k}^{m}\right\}_{k=1}^{\infty}$ where

$$
x_{k}^{m}= \begin{cases}1 & \text { if } k=m \\ 0 & \text { otherwise }\end{cases}
$$

Show that the sequence $\left\{x^{m}\right\}_{m=1}^{\infty}$ in $l^{p}$ (cf. Example $\left.8.9(i i i)\right)$ converges to 0 in the weak topology but does not converge in the strong topology.

### 8.6. Weak and Strong Convergence in $L^{p}$

Although it is easily seen that strong convergence implies weak convergence in $L^{p}$, it is shown below that under certain conditions, weak convergence implies strong convergence.

We now apply some of the results of this chapter in the setting of $L^{p}$ spaces. To begin we note that if $1 \leq p<\infty$, then, in view of Example 8.34 (ii), a sequence $\left\{f_{k}\right\}_{k=1}^{\infty}$ in $L^{p}(X, \mathcal{M}, \mu)$ converges weakly to $f \in L^{p}(X, \mathcal{M}, \mu)$ if and only if

$$
\lim _{k \rightarrow \infty} \int_{X} f_{k} g d \mu=\int_{X} f g d \mu
$$

for each $g \in L^{p^{\prime}}(X, \mathcal{M}, \mu)$.
8.58. Theorem. Let $(X, \mathcal{M}, \mu)$ be a measure space and suppose $f$ and $\left\{f_{k}\right\}_{k=1}^{\infty}$ are functions in $L^{p}(X, \mathcal{M}, \mu)$. If $1 \leq p<\infty$, and $\left\|f_{k}-f\right\|_{p} \rightarrow 0$, then $f_{k} \rightarrow f$ weakly in $L^{p}$.

Proof. This is a consequence of the fact that, in any normed linear space, strong convergence implies week convergence (see Exercise 11, Section 8.4). In this theorem, where the normed linear space is $L^{p}$, the result also follows from Hölder's inequality.

If $\left\{f_{k}\right\}_{k=1}^{\infty}$ is a sequence of functions with $\left\|f_{k}\right\|_{p} \leq M$ for some $M$ and all $k$, then since $L^{p}(X)$ is reflexive for $1<p<\infty$, Theorem 8.37 asserts that there is a subsequence that converges weakly to some $f \in L^{p}(X)$. The next result shows that if it is also known that $f_{k} \rightarrow f \mu$-a.e., then the full sequence converges weakly to $f$.
8.59. Theorem. Let $1<p<\infty$. If $f_{k} \rightarrow f \mu$-a.e., then $f_{k} \rightarrow f$ weakly in $L^{p}$ if and only if $\left\{\left\|f_{k}\right\|_{p}\right\}$ is a bounded sequence.

Proof. Necessity follows immediately from Theorem 8.35.
To prove sufficiency, let $M \geq 0$ be such that $\left\|f_{k}\right\|_{p} \leq M$ for all positive integers $k$. Then Fatou's Lemma implies

$$
\begin{align*}
\|f\|_{p}^{p} & =\int_{X}|f|^{p} \quad d \mu=\int_{X} \lim _{k \rightarrow \infty}\left|f_{k}\right|^{p} d \mu \\
& \leq \liminf _{k \rightarrow \infty} \int_{X}\left|f_{k}\right|^{p} d \mu \leq M^{p} \tag{8.19}
\end{align*}
$$

Let $\varepsilon>0$ and $g \in L^{p^{\prime}}(X)$. Refer to Theorem 6.40 to obtain $\delta>0$ such that

$$
\begin{equation*}
\left(\int_{E}|g|^{p^{\prime}} d \mu\right)^{1 / p^{\prime}}<\frac{\varepsilon}{6 M} \tag{8.20}
\end{equation*}
$$

whenever $E \in \mathcal{M}$ and $\mu(E)<\delta$. We claim there exists a set $F \in \mathcal{M}$ such that $\mu(F)<\infty$ and

$$
\begin{equation*}
\left(\int_{\widetilde{F}}|g|^{p^{\prime}} d \mu\right)^{1 / p^{\prime}}<\frac{\varepsilon}{6 M} . \tag{8.21}
\end{equation*}
$$

To verify the claim, set $A_{t}:=\left\{x:|g(x)|^{p^{\prime}} \geq t\right\}$ and observe that by the Monotone Convergence Theorem,

$$
\lim _{t \rightarrow 0^{+}} \int_{X} \chi_{A_{t}}|g|^{p^{\prime}} d \mu=\int_{X}|g|^{p^{\prime}} d \mu
$$

and thus (8.21) holds with $F=A_{t}$ for sufficiently small positive $t$.
We can now apply Egoroff's Theorem (Theorem 5.18) on $F$ to obtain $A \in \mathcal{M}$ such that $A \subset F, \mu(F-A)<\delta$, and $f_{k} \rightarrow f$ uniformly on $A$. Let $k_{0}$ be such that $k \geq k_{0}$ implies

$$
\begin{equation*}
\left(\int_{A}\left|f-f_{k}\right|^{p} d \mu\right)^{1 / p}\|g\|_{p^{\prime}}<\frac{\varepsilon}{3} \tag{8.22}
\end{equation*}
$$

Setting $E=F-A$ in (8.20), we obtain from (8.19), (8.21), (8.22), and Hölder's inequality,

$$
\begin{aligned}
\left|\int_{X} f g d \mu-\int_{X} f_{k} g d \mu\right| \leq & \int_{X}\left|f-f_{k}\right||g| d \mu \\
= & \int_{A}\left|f-f_{k}\right||g| d \mu+\int_{F-A}\left|f-f_{k}\right||g| d \mu \\
& +\int_{\widetilde{F}}\left|f-f_{k}\right||g| d \mu \\
\leq & \left\|f-f_{k}\right\|_{p, A}\|g\|_{p^{\prime}} \\
& +\left\|f-f_{k}\right\|_{p}\left(\|g\|_{p^{\prime}, F-A}+\|g\|_{p^{\prime}, \widetilde{F}}\right) \\
\leq & \frac{\varepsilon}{3}+2 M\left(\frac{\varepsilon}{6 M}+\frac{\varepsilon}{6 M}\right) \\
= & \varepsilon
\end{aligned}
$$

for all $k \geq k_{0}$.
If the hypotheses of the last result are changed to include that $\left\|f_{k}\right\|_{p} \rightarrow\|f\|_{p}$, then we can prove that $\left\|f_{k}-f\right\|_{p} \rightarrow 0$. This is an immediate consequence of the following theorem.
8.60. Theorem. Let $1 \leq p<\infty$. Suppose $f$ and $\left\{f_{k}\right\}_{k=1}^{\infty}$ are functions in $L^{p}(X, \mathcal{M}, \mu)$ such that $f_{k} \rightarrow f \mu$-a.e. and the sequence $\left\{\left\|f_{k}\right\|_{p}\right\}_{k=1}^{\infty}$ is bounded. Then

$$
\lim _{k \rightarrow \infty}\left(\left\|f_{k}\right\|_{p}^{p}-\left\|f_{k}-f\right\|_{p}^{p}\right)=\|f\|_{p}^{p}
$$

Proof. Set

$$
M:=\sup _{k \geq 1}\left\|f_{k}\right\|_{p}<\infty
$$

and note that by Fatou's Lemma, $\|f\|_{p} \leq M$.
Fix $\varepsilon>0$ and observe that the function

$$
h_{\varepsilon}(t):=\left||t+1|^{p}-|t|^{p}\right|-\varepsilon|t|^{p}
$$

is continuous on $\mathbb{R}$ and

$$
\lim _{|t| \rightarrow \infty} h_{\varepsilon}(t)=-\infty
$$

Thus there is a constant $C_{\varepsilon}>0$ such that $h_{\varepsilon}(t)<C_{\varepsilon}$ for all $t \in \mathbb{R}$. It follows that

$$
\begin{equation*}
\| a+\left.b\right|^{p}-\left.|a|^{p}|\leq \varepsilon| a\right|^{p}+C_{\varepsilon}|b|^{p} \tag{8.23}
\end{equation*}
$$

for any real numbers $a$ and $b$.
Set

$$
G_{k}^{\varepsilon}:=\left[\left|\left|f_{k}\right|^{p}-\left|f_{k}-f\right|^{p}-|f|^{p}\right|-\varepsilon\left|f_{k}-f\right|^{p}\right]^{+}
$$

and note that $G_{k}^{\varepsilon} \rightarrow 0 \mu$-a.e. as $k \rightarrow \infty$.
Setting $a=f_{k}-f$ and $b=f$ in (8.23) we see that
from which it follows that

$$
G_{k}^{\varepsilon} \leq\left(C_{\varepsilon}+1\right)|f|^{p}
$$

and, by the Dominated Convergence Theorem,

$$
\lim _{k \rightarrow \infty} \int_{X} G_{k}^{\varepsilon} d \mu=0
$$

Since

$$
\left|\left|f_{k}\right|^{p}-\left|f_{k}-f\right|^{p}-|f|^{p}\right| \leq G_{k}^{\varepsilon}+\varepsilon\left|f_{k}-f\right|^{p}
$$

we see that

$$
\left.\limsup _{k \rightarrow \infty} \int_{X}| | f_{k}\right|^{p}-\left|f_{k}-f\right|^{p}-|f|^{p} \mid d \mu \leq \varepsilon(2 M)^{p}
$$

Since $\varepsilon$ is arbitrary the proof is complete.
We have the following
8.61. Corollary. Let $1<p<\infty$. If $f_{k} \rightarrow f \mu$-a.e. and $\left\|f_{k}\right\|_{p} \rightarrow\|f\|_{p}$, then $f_{k} \rightarrow f$ weakly in $L^{p}$ and $\left\|f_{k}-f\right\|_{p} \rightarrow 0$.

The following provides a summary of results.
8.62. Theorem. The following statements hold in a general measure space $(X, \mathcal{M}, \mu)$.
(i) If $f_{k} \rightarrow f \mu$-a.e., then

$$
\|f\|_{p} \leq \liminf _{k \rightarrow \infty}\left\|f_{k}\right\|_{p}, \quad 1 \leq p<\infty
$$

(ii) If $f_{k} \rightarrow f \mu$-a.e. and $\left\{\left\|f_{k}\right\|_{p}\right\}$ is a bounded sequence, then $f_{k} \rightarrow f$ weakly in $L^{p}, 1<p<\infty$.
(iii) If $f_{k} \rightarrow f \mu$-a.e. and there exists a function $g \in L^{p}$ such that for each $k$, $\left|f_{k}\right| \leq g \mu$-a.e., on $X$, then $\left\|f_{k}-f\right\|_{p} \rightarrow 0,1 \leq p<\infty$.
(iv) If $f_{k} \rightarrow f \mu$-a.e. and $\left\|f_{k}\right\|_{p} \rightarrow\|f\|_{p}$, then $\left\|f_{k}-f\right\|_{p} \rightarrow 0,1 \leq p<\infty$.

Proof. Fatou's Lemma implies (i), (ii) follows from Theorem 8.59, (iii) follows from Lebesgue's Dominated Convergence Theorem, and (iv) is a restatement of Theorem 8.60.

We conclude this chapter with another proof of the Radon-Nikodym Theorem, which is based on the Riesz Representation Theorem for Hilbert spaces (Theorem 8.49). Since $L^{2}(X, \mu)$ is a Hilbert space, Theorem 8.49 applied to $H=L^{2}(X, \mu)$ is a particular case of Theorem 6.48, which is the Riesz Representation Theorem for $L^{p}$ spaces. The proof of Theorem 6.48 is based on Theorem 6.43 , which proves the Radon-Nikodym Theorem. We note that the shorter proof of the Radon-Nikodym Theorem that we now present does not rely on Theorem 6.48.
8.63. Theorem. Let $\mu$ and $\nu$ be $\sigma$-finite measures on $(X, \mathcal{M})$ with $\nu \ll \mu$. Then there exists a function $h \in L_{\mathrm{loc}}^{1}(X, \mathcal{M}, \mu)$ such that

$$
\nu(E)=\int_{X} h d \mu
$$

for all $E \in \mathcal{M}$.

Proof. First, we will assume $\nu \geq 0$ and that both measures are finite. We will prove that there is a measurable function $g$ on $X$ such that $0 \leq g<1$ and

$$
\int_{X} f(1-g) d \nu=\int_{X} f g d \mu
$$

for all $f \in L^{2}(X, \mathcal{M}, \mu+\nu)$. For this purpose, define

$$
\begin{equation*}
T(f)=\int_{X} f d \nu \tag{8.24}
\end{equation*}
$$

for $f \in L^{2}(X, \mathcal{M}, \mu+\nu)$. By Hölder's inequality, it follows that

$$
f \in L^{1}(X, \mathcal{M}, \mu+\nu)
$$

and therefore $f \in L^{1}(X, \mathcal{M}, \nu)$ also. Hence, we see that $T$ is finite for $f \in$ $L^{2}(X, \mathcal{M}, \mu+\nu)$ and thus is a well-defined linear functional. Furthermore, it is a bounded linear functional because

$$
\begin{aligned}
|T(f)| & \leq\left(\int_{X}|f|^{2} d \nu\right)^{1 / 2}(\nu(X))^{1 / 2} \\
& =\|f\|_{2 ; \nu}(\nu(X))^{1 / 2} \\
& \leq\|f\|_{2 ; \nu+\mu}(\nu(X))^{1 / 2}
\end{aligned}
$$

Referring to Theorem 8.49, there is a function $\varphi \in L^{2}(\mu+\nu)$ such that

$$
\begin{equation*}
T(f)=\int_{X} f \varphi d(\nu+\mu) \tag{8.25}
\end{equation*}
$$

for all $f \in L^{2}(\mu+\nu)$. Observe that $\varphi \geq 0(\mu+\nu)$-a.e, for otherwise we would obtain

$$
T\left(\chi_{A}\right)<0
$$

where $A:=\{\varphi<0\}$, which is impossible. Now, (8.24) and (8.25) imply

$$
\begin{equation*}
\int_{X} f(1-\varphi) d \nu=\int_{X} f \varphi d \mu \tag{8.26}
\end{equation*}
$$

for $f \in L^{2}(\mu+\nu)$. If $f$ is taken as $\chi_{E}$ where $E:=\{\varphi \geq 1\}$, we obtain

$$
0 \leq \mu(E)=\int_{X} \chi_{E} d \mu \leq \int_{X} \chi_{E} \varphi d \mu=\int_{X} \chi_{E}(1-\varphi) d \nu \leq 0
$$

Hence, we have $\mu(E)=0$ and consequently, $\nu(E)=0$. Setting $g=\varphi \chi_{\widetilde{E}}$, we have $0 \leq g<1$ and $g=\varphi$ almost everywhere with respect to both $\mu$ and $\nu$. Reference to (8.26) yields

$$
\int_{X} f(1-g) d \nu=\int_{X} f g d \mu
$$

Since $g$ is bounded, we can replace $f$ by $\left(1+g+g^{2}+\cdots+g^{k}\right) \chi_{E}$ in this equation for any positive integer $k$ and $E$ measurable. Then we obtain

$$
\int_{E}\left(1-g^{k+1}\right) d \nu=\int_{E} g\left(1+g+g^{2}+\cdots+g^{k}\right) d \mu .
$$

Since $0 \leq g<1$ almost everywhere with respect to both $\mu$ and $\nu$, the left side tends to $\nu(E)$ while the integrands on the right side increase monotonically to some function measurable $h$. Thus, by the Monotone Convergence Theorem, we obtain

$$
\nu(E)=\int_{E} h d \mu
$$

for every measurable set $E$. This gives us the desired result in case both $\mu$ and $\nu$ are finite measures.

The proof of the general case proceeds as in Theorem 6.43.

## Exercises for Section 8.6

1. Suppose $f_{i} \rightarrow f$ weakly in $L^{p}(X, \mathcal{M}, \mu), 1<p<\infty$ and that $f_{i} \rightarrow f$ pointwise $\mu$-a.e. Prove that $f_{i}^{+} \rightarrow f^{+}$and $f_{i}^{-} \rightarrow f^{-}$weakly in $L^{p}$.
2. Show that the previous Exercise 1 is false if the hypothesis of pointwise convergence is dropped.
3. Let $C:=C[0,1]$ denote the space of continuous functions on $[0,1]$ endowed with the usual sup norm and let $X$ be a linear subspace of $C$ that is closed relative to the $L^{2}$ norm.
(i) Prove that $X$ is also closed in $C$.
(ii) Show that $\|f\|_{2} \leq\|f\|_{\infty}$ for each $f \in X$
(iii) Prove that there exists $M>0$ such that $\|f\|_{\infty} \leq M\|f\|_{2}$ for all $f \in X$.
(iv) For each $t \in[0,1]$, show that there is a function $g_{t} \in L^{2}$ such that

$$
f(t)=\int g_{t}(x) f(x) d x
$$

for all $f \in X$.
(v) Show that if $f_{k} \rightarrow f$ weakly in $L^{2}$ where $f_{k} \in X$, then $f_{k}(x) \rightarrow f(x)$ for each $x \in[0,1]$.
(vi) As a consequence of the above, show that if $f_{k} \rightarrow f$ weakly in $L^{2}$ where $f_{k} \in X$, then $f_{k} \rightarrow f$ strongly in $L^{2}$; that is $\left\|f_{k}-f\right\|_{2} \rightarrow 0$.
4. A subset $K \subset X$ in a linear space is called convex if for any two points $x, y \in K$, the line segment $t x+(1-t) y, t \in[0,1]$, also belong to $K$.
(i) Prove that the unit ball in any normed linear space is convex.
(ii) If $K$ is a convex set, a point $x_{0} \in K$ is called an extreme point of $K$ if $x_{0}$ is not in the interior of any line segment that lies in $K$; that is, if $x_{0}=t y+(1-t) z$ where $0<t<1$, then either $y$ or $z$ is not in $K$. Show that any $f \in L^{p}[0,1], 1<p<\infty$, with $\|f\|_{p}=1$ is an extreme point of the unit ball.
(iii) Show that the extreme points of the unit ball in $L^{\infty}[0,1]$ are those functions $f$ with $|f(x)|=1$ for a.e. $x$.
(iv) Show that the unit ball in $L^{1}[0,1]$ has no extreme points.

## CHAPTER 9

## Measures and Linear Functionals

### 9.1. The Daniell Integral


#### Abstract

Theorem 6.48 states that a function in $L^{p^{\prime}}$ can be regarded as a bounded linear functional on $L^{p}$. Here we show that a large class of measures can be represented as bounded linear functionals on the space of continuous functions. This is a very important result that has many useful applications and provides a fundamental connection between measure theory and functional analysis.


Suppose $(X, \mathcal{M}, \mu)$ is a measure space. Integration defines an operation that is linear, order preserving, and continuous relative to increasing convergence. Specifically, we have
(i) $\int(k f) d \mu=k \int f d \mu$ whenever $k \in \mathbb{R}$ and $f$ is an integrable function
(ii) $\int(f+g) d \mu=\int f d \mu+\int g d \mu$ whenever $f, g$ are integrable
(iii) $\int f d \mu \leq \int g d \mu$ whenever $f, g$ are integrable functions with $f \leq g \mu$-a.e.
(iv) If $\left\{f_{i}\right\}$ is an nondecreasing sequence of integrable functions, then

$$
\lim _{i \rightarrow \infty} \int f_{i} d \mu=\int \lim _{i \rightarrow \infty} f_{i} d \mu
$$

The main objective of this section is to show that if a linear functional, defined on an appropriate space of functions, possesses the four properties above, then it can be expressed as the operation of integration with respect to some measure. Thus, we will have shown that these properties completely characterize the operation of integration.

It turns out that the proof of our main result is no more difficult when cast in a very general framework, so we proceed by introducing the concept of a lattice.
9.1. Definition. If $f, g$ are real-valued functions defined on a space $X$, we define

$$
\begin{aligned}
(f \wedge g)(x) & :=\min [f(x), g(x)] \\
(f \vee g)(x) & :=\max [f(x), g(x)]
\end{aligned}
$$

A collection $L$ of real-valued functions defined on an abstract space $X$ is called a lattice provided the following conditions are satisfied: if $0 \leq c<\infty$ and $f$ and $g$
are elements of $L$, then so are the functions $f+g, c f, f \wedge g$, and $f \wedge c$. Furthermore, if $f \leq g$, then $g-f$ is required to belong to $L$. Note that if $f$ and $g$ belong to $L$ with $g \geq 0$, then so does $f \vee g=f+g-f \wedge g$. Therefore, $f^{+}$belongs to $L$. We define $f^{-}=f^{+}-f$ and so $f^{-} \in L$. We let $L^{+}$denote those functions $f$ in $L$ for which $f \geq 0$. Clearly, if $L$ is a lattice, then so is $L^{+}$.

For example, the space of continuous functions on a metric space is a lattice as well as the space of integrable functions, but the collection of lower semicontinuous functions is not a lattice because it is not closed under the $\wedge$-operation (see Exercise 9.1).
9.2. Theorem. Suppose $L$ is a lattice of functions on $X$ and let $T: L \rightarrow \mathbb{R}$ be a functional satisfying the following conditions for all functions in $L$ :
(i) $T(f+g)=T(f)+T(g)$
(ii) $T(c f)=c T(f)$ whenever $0 \leq c<\infty$
(iii) $T(f) \geq T(g)$ whenever $f \geq g$
(iv) $T(f)=\lim _{i \rightarrow \infty} T\left(f_{i}\right)$ whenever $f:=\lim _{i \rightarrow \infty} f_{i}$ is a member of $L$
(v) and $\left\{f_{i}\right\}$ is nondecreasing.

Then, there exists an outer measure $\mu$ on $X$ such that for each $f \in L, f$ is $\mu$ measurable and

$$
T(f)=\int_{X} f d \mu
$$

In particular, $\{f>t\}$ is $\mu$-measurable whenever $f \in L$ and $t \in \mathbb{R}$.
Proof. First, observe that since $T(0)=T(0 \cdot f)=0 \cdot T(f)=0$, (iii) above implies $T(f) \geq 0$ whenever $f \in L^{+}$. Next, with the convention that the infimum of the empty set is $\infty$, for an arbitrary set $A \subset X$, we define

$$
\mu(A)=\inf \left\{\lim _{i \rightarrow \infty} T\left(f_{i}\right)\right\}
$$

where the infimum is taken over all sequences of functions $\left\{f_{i}\right\}$ with the property

$$
\begin{equation*}
f_{i} \in L^{+},\left\{f_{i}\right\}_{i=1}^{\infty} \text { is nondecreasing, and } \lim _{i \rightarrow \infty} f_{i} \geq \chi_{A} \tag{9.1}
\end{equation*}
$$

Such sequences of functions are called admissible for $A$. In accordance with our convention, if there is no admissible sequence of functions for $A$, we define $\mu(A)=$ $\infty$. It follows from definition that if $f \in L^{+}$and $f \geq \chi_{A}$, then

$$
\begin{equation*}
T(f) \geq \mu(A) \tag{9.2}
\end{equation*}
$$

On the other hand, if $f \in L^{+}$and $f \leq \chi_{A}$, then

$$
\begin{equation*}
T(f) \leq \mu(A) \tag{9.3}
\end{equation*}
$$

This is true because if $\left\{f_{i}\right\}$ is admissible for $A$, then $g_{i}:=f_{i} \wedge f$ is a nondecreasing sequence with $\lim _{i \rightarrow \infty} g_{i}=f$. Hence,

$$
T(f)=\lim _{i \rightarrow \infty} T\left(g_{i}\right) \leq \lim _{i \rightarrow \infty} T\left(f_{i}\right)
$$

and therefore

$$
T(f) \leq \mu(A)
$$

The first step is to show that $\mu$ is an outer measure on $X$. For this, the only nontrivial property to be established is countable subadditivity. For this purpose, let

$$
A \subset \bigcup_{i=1}^{\infty} A_{i} .
$$

For each fixed $i$ and arbitrary $\varepsilon>0$, let $\left\{f_{i, j}\right\}_{j=1}^{\infty}$ be an admissible sequence of functions for $A_{i}$ with the property that

$$
\lim _{j \rightarrow \infty} T\left(f_{i, j}\right)<\mu\left(A_{i}\right)+\frac{\varepsilon}{2^{i}} .
$$

Now define

$$
g_{k}=\sum_{i=1}^{k} f_{i, k}
$$

and obtain

$$
\begin{aligned}
T\left(g_{k}\right) & =\sum_{i=1}^{k} T\left(f_{i, k}\right) \\
& \leq \sum_{i=1}^{k} \lim _{m \rightarrow \infty} T\left(f_{i, m}\right) \\
& \leq \sum_{i=1}^{k}\left(\mu\left(A_{i}\right)+\frac{\varepsilon}{2^{i}}\right) .
\end{aligned}
$$

After we show that $\left\{g_{k}\right\}$ is admissible for $A$, we can take limits of both sides as $k \rightarrow \infty$ and conclude

$$
\mu(A) \leq \sum_{i=1}^{\infty} \mu\left(A_{i}\right)+\varepsilon
$$

Since $\varepsilon$ is arbitrary, this would prove the countable subadditivity of $\mu$ and thus establish that $\mu$ is an outer measure on $X$.

To see that $\left\{g_{k}\right\}$ is admissible for $A$, select $x \in A$. Then $x \in A_{i}$ for some $i$ and with $k:=\max (i, j)$, we have $g_{k}(x) \geq f_{i, j}(x)$. Hence,

$$
\lim _{k \rightarrow \infty} g_{k}(x) \geq \lim _{j \rightarrow \infty} f_{i, j}(x) \geq 1
$$

The next step is to prove that each element $f \in L$ is a $\mu$-measurable function. Since $f=f^{+}-f^{-}$, it suffices to show that $f^{+}$is $\mu$-measurable, the proof involving
$f^{-}$being similar. For this we appeal to Theorem 5.13 which asserts that it is sufficient to prove

$$
\begin{equation*}
\mu(A) \geq \mu\left(A \cap\left\{f^{+} \leq a\right\}\right)+\mu\left(A \cap\left\{f^{+} \geq b\right\}\right) \tag{9.4}
\end{equation*}
$$

whenever $A \subset X$ and $a<b$ are real numbers. Since $f^{+} \geq 0$, we may as well take $a \geq 0$. Let $\left\{g_{i}\right\}$ be an admissible sequence of functions for $A$ and define

$$
h=\frac{\left[f^{+} \wedge b-f^{+} \wedge a\right]}{b-a}, \quad k_{i}=g_{i} \wedge h
$$

Observe that $h=1$ on $\left\{f^{+} \geq b\right\}$. Since $\left\{g_{i}\right\}$ is admissible for $A$ it follows that $\left\{k_{i}\right\}$ is admissible for $A \cap\left\{f^{+} \geq b\right\}$. Furthermore, since $h=0$ on $\left\{f^{+} \leq a\right\}$, we have that $\left\{g_{i}-k_{i}\right\}$ is admissible for $A \cap\left\{f^{+} \leq a\right\}$. Consequently,

$$
\lim _{i \rightarrow \infty} T\left(g_{i}\right)=\lim _{i \rightarrow \infty}\left[T\left(k_{i}\right)+T\left(g_{i}-k_{i}\right)\right] \geq \mu\left(A \cap\left\{f^{+} \geq b\right\}\right)+\mu\left(A \cap\left\{f^{+} \leq a\right\}\right)
$$

Since $\left\{g_{i}\right\}$ is an arbitrary admissible sequence for $A$, we obtain

$$
\mu(A) \geq \mu\left(A \cap\left\{f^{+} \geq b\right\}\right)+\mu\left(A \cap\left\{f^{+} \leq a\right\}\right)
$$

which proves (9.4).
The last step is to prove that

$$
T(f)=\int f d \mu \quad \text { for } \quad f \in L
$$

We begin by considering $f \in L^{+}$. With $f_{t}:=f \wedge t, t \geq 0$, note that if $\varepsilon>0$ and $k$ is a positive integer, then

$$
0 \leq f_{k \varepsilon}(x)-f_{(k-1) \varepsilon}(x) \quad \text { for } \quad x \in X
$$

and

$$
f_{k \varepsilon}(x)-f_{(k-1) \varepsilon}(x)= \begin{cases}\varepsilon & \text { for } f(x) \geq k \varepsilon \\ 0 & \text { for } f(x) \leq(k-1) \varepsilon\end{cases}
$$

Therefore,

$$
\begin{array}{rlrl}
T\left(f_{k \varepsilon}-f_{(k-1) \varepsilon}\right) & \geq \varepsilon \mu(\{f \geq k \varepsilon\}) & & \text { by }(9.2) \\
& \geq \int\left(f_{(k+1) \varepsilon}-f_{k \varepsilon}\right) d \mu & & \text { since } f_{(k+1) \varepsilon}-f_{k \varepsilon} \leq \varepsilon \chi_{\{f \geq k \varepsilon\}} \\
& \geq \varepsilon \mu(\{f \geq(k+1) \varepsilon\}) & & \text { since } f_{(k+1) \varepsilon}-f_{k \varepsilon}=\varepsilon \chi_{\{f \geq k \varepsilon\}} \\
& & \text { on }\{f \geq(k+1) \varepsilon\} \\
& \geq T\left(f_{(k+2) \varepsilon}-f_{(k+1) \varepsilon}\right) . & & \text { by }(9.3)
\end{array}
$$

Now taking the sum as $k$ ranges from 1 to $n$, we obtain

$$
T\left(f_{n \varepsilon}\right) \geq \int\left(f_{(n+1) \varepsilon}-f_{\varepsilon}\right) d \mu \geq T\left(f_{(n+2) \varepsilon}-f_{2 \varepsilon}\right)
$$

Since $\left\{f_{n \varepsilon}\right\}$ is a nondecreasing sequence with $\lim _{n \rightarrow \infty} f_{n \varepsilon}=f$, we have

$$
T(f) \geq \int\left(f-f_{\varepsilon}\right) d \mu \geq T\left(f-f_{2 \varepsilon}\right)
$$

Also, $f_{\varepsilon} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$so that

$$
T(f)=\int f d \mu
$$

Finally, if $f \in L$, then $f^{+} \in L^{+}, f^{-} \in L^{+}$, thus yielding

$$
T(f)=T\left(f^{+}\right)-T\left(f^{-}\right)=\int f^{+} d \mu-\int f^{-} d \mu=\int f d \mu
$$

Now that we have established the existence of an outer measure $\mu$ corresponding to the functional $T$, we address the question of its uniqueness. For this purpose, we need the following lemma, which asserts that $\mu$ possesses a type of outer regularity.
9.3. Lemma. Under the assumptions of the preceding Theorem, let $\mathcal{M}$ denote the $\sigma$-algebra generated by all sets of the form $\{f>t\}$ where $f \in L$ and $t \in \mathbb{R}$. Then, for any $A \subset X$ there is a set $W \in \mathcal{M}$ such that

$$
A \subset W \quad \text { and } \quad \mu(A)=\mu(W)
$$

Proof. If $\mu(A)=\infty$, take $W=X$. If $\mu(A)<\infty$, we proceed as follows. For each positive integer $i$ let $\left\{f_{i, j}\right\}_{j=1}^{\infty}$ be an admissible sequence for $A$ with the property

$$
\lim _{j \rightarrow \infty} T\left(f_{i, j}\right)<\mu(A)+\frac{1}{i}
$$

and define, for each positive integer $j$,

$$
\begin{aligned}
g_{i, j}: & =\inf \left\{f_{1, j}, f_{2, j}, \ldots, f_{i, j}\right\} \\
B_{i, j}: & =\left\{x: g_{i, j}(x)>1-\frac{1}{i}\right\}
\end{aligned}
$$

Then,

$$
g_{i+1, j} \leq g_{i, j} \leq g_{i, j+1} \quad \text { and } \quad B_{i+1, j} \subset B_{i, j} \subset B_{i, j+1}
$$

Furthermore, with the help of (9.2),

$$
(1-1 / i) \mu\left(B_{i, j}\right) \leq T\left(g_{i, j}\right) \leq T\left(f_{i, j}\right)
$$

Now let

$$
V_{i}=\bigcup_{j=1}^{\infty} B_{i, j} .
$$

Observe that $V_{i} \supset V_{i+1} \supset A$ and $V_{i} \in \mathcal{M}$ for all $i$. Indeed, to verify that $V_{i} \supset A$, it is sufficient to show $g_{i, j}\left(x_{0}\right)>1-1 / i$ for $x_{0} \in A$ whenever $j$ is sufficiently large. This is accomplished by observing that there exist positive integers $j_{1}, j_{2}, \ldots, j_{i}$ such that $f_{1, j}\left(x_{0}\right)>1-1 / i$ for $j \geq j_{1}, f_{2, j}\left(x_{0}\right)>1-1 / i$ for $j \geq j_{2}, \ldots, f_{i, j}\left(x_{0}\right)>1-1 / i$
for $j \geq j_{i}$. Thus, for $j$ larger than $\max \left\{j_{1}, j_{2}, \ldots, j_{i}\right\}$, we have $g_{i, j}\left(x_{0}\right)>1-1 / i$. For each positive integer $i$,

$$
\begin{align*}
\mu(A) \leq \mu\left(V_{i}\right) & =\lim _{j \rightarrow \infty} \mu\left(B_{i, j}\right) \\
& \leq(1-1 / i)^{-1} \lim _{j \rightarrow \infty} T\left(f_{i, j}\right)  \tag{9.5}\\
& \leq(1-1 / i)^{-1}[\mu(A)+1 / i]
\end{align*}
$$

and therefore,

$$
\begin{equation*}
\mu(A)=\lim _{i \rightarrow \infty} \mu\left(V_{i}\right) \tag{9.6}
\end{equation*}
$$

Note that (9.5) implies $\mu\left(V_{i}\right)<\infty$. Therefore, if we take

$$
W=\bigcap_{i=1}^{\infty} V_{i}
$$

we obtain $A \subset W$, $W \in \mathcal{M}$, and

$$
\begin{aligned}
\mu(W) & =\mu\left(\bigcap_{i=1}^{\infty} V_{i}\right) \\
& =\lim _{i \rightarrow \infty} \mu\left(V_{i}\right) \\
& =\mu(A)
\end{aligned}
$$

The preceding proof reveals the manner in which $T$ uniquely determines $\mu$. Indeed, for $f \in L^{+}$and for $t, h>0$, define

$$
\begin{equation*}
f_{h}=\frac{f \wedge(t+h)-f \wedge t}{h} \tag{9.7}
\end{equation*}
$$

Let $\left\{h_{i}\right\}$ be a nonincreasing sequence of positive numbers such that $h_{i} \rightarrow 0$. Observe that $\left\{f_{h_{i}}\right\}$ is admissible for the set $\{f>t\}$ while $f_{h_{i}} \leq \chi_{\{f>t\}}$ for all small $h_{i}$. From Theorem 9.2 we know that

$$
T\left(f_{h_{i}}\right)=\int f_{h_{i}} d \mu
$$

Thus, the Monotone Convergence Theorem implies that

$$
\lim _{i \rightarrow \infty} T\left(f_{h_{i}}\right)=\mu(\{f>t\})
$$

and this shows that the value of $\mu$ on the set $\{f>t\}$ is uniquely determined by $T$. From this it follows that $\mu\left(B_{i, j}\right)$ is uniquely determined by $T$, and therefore by (9.5), the same is true for $\mu\left(V_{i}\right)$. Finally, referring to (9.6), where it is assumed that $\mu(A)<\infty$, we have that $\mu(A)$ is uniquely determined by $T$. The requirement that $\mu(A)<\infty$ will be be ensured if $\chi_{A} \leq f$ for some $f \in L$, see (9.2). Thus, we have the following corollary.
9.4. Corollary. If $T$ is as in Theorem 9.2, the corresponding outer measure $\mu$ is uniquely determined by $T$ on all sets $A$ with the property that $\chi_{A} \leq f$ for some $f \in L$.

Functionals that satisfy the conditions of Theorem 9.2 are called monotone (or alternatively, positive). In the spirit of the Jordan Decomposition Theorem, we now investigate the question of determining those functionals that can be written as the difference of monotone functionals.
9.5. Theorem. Suppose $L$ is a lattice of functions on $X$ and let $T: L \rightarrow \mathbb{R}$ be a functional satisfying the following four conditions for all functions in $L$ :

$$
\begin{aligned}
T(f+g) & =T(f)+T(g) \\
T(c f) & =c T(f) \text { whenever } 0 \leq c<\infty \\
\sup \{T(k) & : 0 \leq k \leq f\}<\infty \text { for all } f \in L^{+} \\
T(f) & =\lim _{i \rightarrow \infty} T\left(f_{i}\right) \text { whenever } f=\lim _{i \rightarrow \infty} f_{i} \text { and }\left\{f_{i}\right\} \text { is nondecreasing. }
\end{aligned}
$$

Then, there exist positive functionals $T^{+}$and $T^{-}$defined on the lattice $L^{+}$satisfying the conditions of Theorem 9.2 and the property

$$
\begin{equation*}
T(f)=T^{+}(f)-T^{-}(f) \tag{9.8}
\end{equation*}
$$

for all $f \in L^{+}$.
Proof. We define $T^{+}$and $T^{-}$on $L^{+}$as follows:

$$
\begin{aligned}
& T^{+}(f)=\sup \{T(k): 0 \leq k \leq f\} \\
& T^{-}(f)=-\inf \{T(k): 0 \leq k \leq f\}
\end{aligned}
$$

To prove (9.8), let $f, g \in L^{+}$with $f \geq g$. Then $f \geq f-g$ and therefore $-T^{-}(f) \leq$ $T(f-g)$. Hence, $T(g)-T^{-}(f) \leq T(g)+T(f-g)=T(f)$. Taking the supremum over all $g$ with $g \leq f$ yields

$$
T^{+}(f)-T^{-}(f) \leq T(f)
$$

Similarly, since $f-g \leq f, T(f-g) \leq T^{+}(f)$ so that $T(f)=T(g)+T(f-g) \leq$ $T(g)+T^{+}(f)$. Taking the infimum over all $0 \leq g \leq f$ implies

$$
T(f) \leq-T^{-}(f)+T^{+}(f)
$$

Hence we obtain

$$
T(f)=T^{+}(f)-T^{-}(f)
$$

Now we will prove that $T^{+}$satisfies the conditions of Theorem 9.2 on $L^{+}$. For this, let $f, g, h \in L^{+}$with $f+g \geq h$ and set $k=\inf \{f, h\}$. Then $f \geq k$ and $g \geq h-k$; consequently,

$$
T^{+}(f)+T^{+}(g) \geq T(k)+T(h-k)=T(h)
$$

Since this holds for all $h \leq f+g$ with $h \in L^{+}$, we have

$$
T^{+}(f)+T^{+}(g) \geq T^{+}(f+g)
$$

It is easy to verify the opposite inequality and also that $T^{+}$is both positively homogeneous and monotone.

To show that $T^{+}$satisfies the last condition of Theorem 9.2 , let $\left\{f_{i}\right\}$ be a nondecreasing sequence in $L^{+}$such that $f_{i} \rightarrow f$. If $k \in L^{+}$and $k \leq f$, then the sequence $g_{i}=\inf \left\{f_{i}, k\right\}$ is nondecreasing and converges to $k$ as $i \rightarrow \infty$. Hence

$$
T(k)=\lim _{i \rightarrow \infty} T\left(g_{i}\right) \leq \lim _{i \rightarrow \infty} T^{+}\left(f_{i}\right)
$$

and therefore

$$
T^{+}(f) \leq \lim _{i \rightarrow \infty} T^{+}\left(f_{i}\right)
$$

The opposite inequality is obvious and thus equality holds.
That $T^{-}$also satisfies the conditions of Theorem 9.2 is almost immediate. Indeed, let $R=-T$ and observe that $R$ satisfies the first, second, and fourth conditions of our current theorem. It also satisfies the third because $R^{+}(f)=$ $(-T)^{+}(f)=T^{-}(f)=T(f)-T^{+}(f)<\infty$ for each $f \in L^{+}$. Thus what we have just proved for $T^{+}$applies to $R^{+}$as well. In particular, $R^{+}$satisfies the conditions of Theorem 9.2.
9.6. Corollary. Let $T$ be as in the previous theorem. Then there exist outer measures $\mu^{+}$and $\mu^{-}$on $X$ such that for each $f \in L, f$ is both $\mu^{+}$and $\mu^{-}$measurable and

$$
T(f)=\int f d \mu^{+}-\int f d \mu^{-}
$$

Proof. Theorem 9.2 supplies outer measures $\mu^{+}, \mu^{-}$such that

$$
\begin{aligned}
& T^{+}(f)=\int f d \mu^{+} \\
& T^{-}(f)=\int f d \mu^{-}
\end{aligned}
$$

for all $f \in L^{+}$. Since each $f \in L$ can be written as $f=f^{+}-f^{-}$, the result follows.

A functional $T$ satisfying the conditions of Theorem 9.2 is called a monotone Daniell integral. By a Daniell integral, we mean a functional $T$ of the type in Theorem 9.5.

## Exercises for Section 9.1

1. Show that the space of lower semicontinuous functions on a metric space is not a lattice because it is not closed under the $\wedge$-operation.
2. Let $L$ denote the family of all functions of the form $u \circ p$ where $u: \mathbb{R} \rightarrow \mathbb{R}$ is continuous and $p: \mathbb{R}^{2} \rightarrow \mathbb{R}$ is the orthogonal projection defined by $p\left(x_{1}, x_{2}\right)=x_{1}$. Define a functional $T$ on $L$ by

$$
T(u \circ p)=\int_{0}^{1} u d \lambda
$$

Prove that $T$ meets the conditions of Theorem 9.2 and that the measure $\mu$ representing $T$ satisfies

$$
\mu(A)=\lambda[p(A) \cap[0,1]]
$$

for all $A \subset \mathbb{R}^{2}$. Also, show that not all Borel sets are $\mu$-measurable.
3. Prove that the sequence $f_{h_{i}}$ is admissible for the set $\{f>t\}$ where $f_{h}$ is defined by (9.7).
4. Let $L=C_{c}(\mathbb{R})$.
(i) Prove Dini's Theorem: If $\left\{f_{i}\right\}$ is a nonincreasing sequence in $L$ that converges pointwise to 0 , then $f_{i} \rightarrow 0$ uniformly.
(ii) Define $T: L \rightarrow \mathbb{R}$ by

$$
T(f)=\int f
$$

where the integral denotes the Riemann integral. Prove that $T$ is a Daniell integral.
5. Roughly speaking, it can be shown that the measures $\mu^{+}$and $\mu^{-}$that occur in Corollary 9.6 are carried by disjoint sets. This can be made precise in the following way. Prove, for an arbitrary $f \in L^{+}$, there exists a $\mu$-measurable function $g$ on $X$ such that

$$
g(x)= \begin{cases}f(x) & \text { for } \mu^{+} \text {-a.e. } x \\ 0 & \text { for } \mu^{-} \text {-a.e. } x\end{cases}
$$

### 9.2. The Riesz Representation Theorem


#### Abstract

As a major application of the theorem in the previous section, we will prove the Riesz Representation theorem which asserts that a large class of measures on a locally compact Hausdorff space can be identified with positive linear functionals on the space of continuous functions with compact support.


We recall that a Hausdorff space $X$ is said to be locally compact if for each $x \in X$, there is an open set $U$ containing $x$ such that $\bar{U}$ is compact. We let

$$
C_{c}(X)
$$

denote the set of all continuous maps $f: X \rightarrow \mathbb{R}$ whose support

$$
\operatorname{spt} f:=\operatorname{closure}\{x: f(x) \neq 0\}
$$

is compact. The class of Baire sets is defined as the smallest $\sigma$-algebra containing the sets $\{f>t\}$ for all $f \in C_{c}(X)$ and all real numbers $t$. Since each $\{f>t\}$ is an open set, it follows that the Borel sets contain the Baire sets. In case $X$ is a locally compact Hausdorff space satisfying the second axiom of countability, the converse is true (Exercise 3, Section 9.2). We will call an outer measure $\mu$ on $X$ a Baire outer measure if all Baire sets are $\mu$-measurable.

We first establish two results that are needed for our development.
9.7. Lemma (Urysohn's Lemma). Let $K \subset U$ be compact and open sets, respectively, in a locally compact Hausdorff space $X$. Then there exists an $f \in C_{c}(X)$ such that

$$
\chi_{K} \leq f \leq \chi_{U}
$$

Proof. Let $r_{1}, r_{2}, \ldots$ be an enumeration of the rationals in $(0,1)$ where we take $r_{1}=0$ and $r_{2}=1$. By Theorem 3.18 , there exist open sets $V_{0}$ and $V_{1}$ with compact closures such that

$$
K \subset V_{1} \subset \bar{V}_{1} \subset V_{0} \subset \bar{V}_{0} \subset V
$$

Proceeding by induction, we will associate an open set $V_{r_{k}}$ with each rational $r_{k}$ in the following way. Assume that $V_{r_{1}}, \ldots, V_{r_{k}}$ have been chosen in such a manner that $\bar{V}_{r_{j}} \subset V_{r_{i}}$ whenever $r_{i}<r_{j}$. Among the numbers $r_{1}, r_{2}, \ldots, r_{k}$ let $r_{m}$ denote the smallest that is larger than $r_{k+1}$ and let $r_{n}$ denote the largest that is smaller than $r_{k+1}$. Referring again to Theorem 3.18, there exists $V_{r_{k+1}}$ such that

$$
\bar{V}_{r_{n}} \subset V_{r_{k+1}} \subset \bar{V}_{r_{k+1}} \subset V_{r_{m}}
$$

Continuing this process, for rational number $r \in[0,1]$ there is a corresponding open set $V_{r}$. The countable collection $\left\{V_{r}\right\}$ satisfies the following properties: $K \subset$ $V_{1}, \bar{V}_{0} \subset U, \bar{V}_{r}$ is compact, and

$$
\begin{equation*}
\bar{V}_{s} \subset V_{r} \quad \text { for } \quad r<s \tag{9.9}
\end{equation*}
$$

For rational numbers $r$ and $s$, define

$$
f_{r}:=r \chi_{V_{r}} \quad \text { and } \quad g_{s}:=\chi_{\bar{V}_{s}}+s \chi_{X-\bar{V}_{s}}
$$

Referring to Definition 3.68, it is easy to see that each $f_{r}$ is lower semicontinuous and each $g_{s}$ is upper semicontinuous. Consequently, with

$$
f(x):=\sup \left\{f_{r}(x): r \in \mathbb{Q} \cap(0,1)\right\} \quad \text { and } \quad g:=\inf \left\{g_{s}(x): r \in \mathbb{Q} \cap(0,1)\right\}
$$

it follows that $f$ is lower semicontinuous and $g$ is upper semicontinuous. Note that $0 \leq f \leq 1, f=1$ on $K$, and that $\operatorname{spt} f \subset \bar{V}_{0}$.

The proof will be completed when we show that $f$ is continuous. This is established by showing $f=g$. To this end, note that $f \leq g$, for otherwise there exists $x \in X$ such that $f_{r}(x)>g_{s}(x)$ for some rationals $r$ and $s$. But this is possible only if $r>s, x \in V_{r}$, and $x \notin \bar{V}_{s}$. However, $r>s$ implies $V_{r} \subset V_{s}$.

Finally, we observe that if $f(x)<g(x)$ for some $x$, then there would exist rationals $r$ and $s$ such that

$$
f(x)<r<s<g(x)
$$

This would imply $x \notin V_{r}$ and $x \in \bar{V}_{s}$, contradicting (9.9). Hence, $f=g$.
9.8. ThEOREM. Suppose $K$ is compact and $V_{1}, \ldots, V_{n}$ are open sets in a locally compact Hausdorff space $X$ such that

$$
K \subset V_{1} \cup \cdots \cup V_{n}
$$

Then there are continuous functions $g_{i}$ with $0 \leq g_{i} \leq 1$ and spt $g_{i} \subset V_{i}, i=$ $1,2, \ldots, n$ such that

$$
g_{1}(x)+g_{2}(x)+\cdots+g_{n}(x)=1 \quad \text { for all } \quad x \in K .
$$

Proof. By Theorem 3.18, each $x \in K$ is contained in an open set $U_{x}$ with compact closure such that $\bar{U}_{x} \subset V_{i}$ for some $i$. Since $K$ is compact, there exist finitely many points $x_{1}, x_{2}, \ldots, x_{k}$ in $K$ such that

$$
K \subset \bigcup_{i=1}^{k} U_{x_{i}}
$$

For $1 \leq j \leq n$, let $G_{j}$ denote the union of those $\bar{U}_{x_{i}}$ that lie in $V_{j}$. Urysohn's Lemma, Lemma 9.7, provides continuous functions $f_{j}$ with $0 \leq f_{j} \leq 1$ such that

$$
\chi_{G_{j}} \leq f_{j} \leq \chi_{V_{j}}
$$

Now $\sum_{j=1}^{n} f_{j} \geq 1$ on $K$, so applying Urysohn's Lemma again, there exists $f \in$ $C_{c}(X)$ with $f=1$ on $K$ and spt $f \subset\left\{\sum_{j=1}^{n} f_{j}>0\right\}$. Let $f_{n+1}=1-f$ so that $\sum_{j=1}^{n+1} f_{j}>0$ everywhere. Now define

$$
g_{i}=\frac{f_{i}}{\sum_{j=1}^{n+1} f_{j}}
$$

to obtain the desired conclusion.

The theorem on the existence of the Daniell integral provides the following as an immediate application.
9.9. Theorem (Riesz Representaton Theorem). Let $X$ be a locally compact Hausdorff space and let $L$ denote the lattice $C_{c}(X)$ of continuous functions with compact support. If $T: L \rightarrow \mathbb{R}$ is a linear functional such that

$$
\begin{equation*}
\sup \{T(g): 0 \leq g \leq f\}<\infty \tag{9.10}
\end{equation*}
$$

whenever $f \in L^{+}$, then there exist Baire outer measures $\mu^{+}$and $\mu^{-}$such that for each $f \in C_{c}(X)$,

$$
T(f)=\int_{X} f d \mu^{+}-\int_{X} f d \mu^{-}
$$

The outer measures $\mu^{+}$and $\mu^{-}$are uniquely determined on the family of compact sets.

Proof. If we can show that $T$ satisfies the hypotheses of Theorem 9.5, then Corollary 9.6 provides outer measures $\mu^{+}$and $\mu^{-}$satisfying our conclusion.

All the conditions of Theorem 9.5 are easily verified except perhaps the last one. Thus, let $\left\{f_{i}\right\}$ be a nondecreasing sequence in $L^{+}$, whose limit is $f \in L^{+}$and refer to Urysohn's Lemma to find a function $g \in C_{c}(X)$ such that $g \geq 1$ on spt $f$.

Choose $\varepsilon>0$ and define compact sets

$$
K_{i}=\left\{x: f(x) \geq f_{i}(x)+\varepsilon\right\},
$$

so that $K_{1} \supset K_{2} \supset \ldots$. Now

$$
\bigcap_{i=1}^{\infty} K_{i}=\emptyset
$$

and therefore the $\widetilde{K}_{i}$ form an open covering of $X$, and in particular, of $\operatorname{spt} f$. Since spt $f$ is compact, there is an index $i_{0}$ such that

$$
\operatorname{spt} f \subset \bigcup_{i=1}^{i_{0}} \widetilde{K_{i}}=\widetilde{K_{i_{0}}}
$$

Since $\left\{f_{i}\right\}$ is a nondecreasing sequence and $g \geq 1$ on spt $f$, this implies that

$$
f(x)<f_{i}(x)+\varepsilon g(x)
$$

for all $i \geq i_{0}$ and all $x \in X$. Note (9.10) implies that

$$
M:=\sup \{|T(k)|: k \in L, 0 \leq k \leq g\}<\infty
$$

Thus, we obtain

$$
0 \leq f-f_{i} \leq \varepsilon g, \quad\left|T\left(f-f_{i}\right)\right| \leq \varepsilon M<\infty
$$

Since $\varepsilon$ is arbitrary, we conclude $T\left(f_{i}\right) \rightarrow T(f)$ as required.

Concerning the assertion of uniqueness, select a compact set $K$ and use Theorem 3.18 to find an open set $U \supset K$ with compact closure. Consider the lattice

$$
L_{U}:=C_{c}(X) \cap\{f: \operatorname{spt} f \subset U\}
$$

and let $T_{U}(f):=T(f)$ for $f \in L_{U}$. Then we obtain outer measures $\mu_{U}^{+}$and $\mu_{U}^{-}$ such that

$$
T_{U}(f)=\int f d \mu_{U}^{+}-\int f d \mu_{U}^{-}
$$

for all $f \in L_{U}$. With the help of Urysohn's Lemma and Corollary 9.4, we find that $\mu_{U}^{+}$and $\mu_{U}^{-}$are uniquely determined and so $\mu_{U}^{+}(K)=\mu^{+}(K)$ and $\mu_{U}^{-}(K)=$ $\mu^{-}(K)$.
9.10. Remark. In the previous result, it might be tempting to define a signed measure $\mu$ by $\mu:=\mu^{+}-\mu^{-}$and thereby reach the conclusion that

$$
T(f)=\int_{X} f d \mu
$$

for each $f \in C_{c}(X)$. However, this is not possible because $\mu^{+}-\mu^{-}$may be undefined. That is, both $\mu^{+}(E)$ and $\mu^{-}(E)$ may possibly assume the value $+\infty$ for some set $E$. See Exercise 6 , Section 9.2 , to obtain a resolution of this in some situations.

If $X$ is a compact Hausdorff space, then all continuous functions are bounded and $C(X)$ becomes a normed linear space with the norm $\|f\|:=\sup \{|f(x)|: x \in$ $X\}$. We will show that there is a very useful characterization of the dual of $C(X)$ that results as a direct consequence of the Riesz Representation Theorem. First, recall that the norm of a linear functional $T$ on $C(X)$ is defined in the usual way by

$$
\|T\|=\sup \{T(f):\|f\| \leq 1\}
$$

If $T$ is written as the difference of two positive functionals, $T=T^{+}{ }^{-} T^{-}$, then

$$
\|T\| \leq\left\|T^{+}\right\|+\left\|T^{-}\right\|=T^{+}(1)+T^{-}(1)
$$

The opposite inequality is also valid, for if $f \in C(X)$ is any function with $0 \leq f \leq 1$, then $|2 f-1| \leq 1$ and

$$
\|T\| \geq T(2 f-1)=2 T(f)-T(1)
$$

Taking the supremum over all such $f$ yields

$$
\begin{aligned}
\|T\| & \geq 2 T^{+}(1)-T(1) \\
& =T^{+}(1)+T^{-}(1)
\end{aligned}
$$

Hence,

$$
\begin{equation*}
\|T\|=T^{+}(1)+T^{-}(1) \tag{9.11}
\end{equation*}
$$

9.11. Corollary. Let $X$ be a compact Hausdorff space. Then for every bounded linear functional $T: C(X) \rightarrow \mathbb{R}$, there exists a unique, signed, Baire outer measure $\mu$ on $X$ such that

$$
T(f)=\int_{X} f d \mu
$$

for each $f \in C(X)$. Moreover, $\|T\|=\|\mu\|(X)$. Thus, the dual of $C(X)$ is isometrically isomorphic to the space of signed, Baire outer measures on $X$.

This will be used in the proof below.
Proof. A bounded linear functional $T$ on $C(X)$ is easily seen to imply property (9.10), and therefore Theorem 9.9 implies there exist Baire outer measures $\mu^{+}$ and $\mu^{-}$such that, with $\mu:=\mu^{+}-\mu^{-}$,

$$
T(f)=\int_{X} f d \mu
$$

for all $f \in C(X)$. Thus,

$$
\begin{aligned}
|T(f)| & \leq \int|f| d\|\mu\| \\
& \leq\|f\|\|\mu\|(X)
\end{aligned}
$$

and therefore, $\|T\| \leq\|\mu\|(X)$. On the other hand, with

$$
T^{+}(f):=\int_{X} f d \mu^{+} \quad \text { and } \quad T^{-}(f):=\int_{X} f d \mu^{-}
$$

we have with the help of (9.11)

$$
\begin{aligned}
\|\mu\|(X) & \leq \mu^{+}(X)+\mu^{-}(X) \\
& =T^{+}(1)+T^{-}(1)=\|T\|
\end{aligned}
$$

Hence, $\|T\|=\|\mu\|(X)$.
Baire sets arise naturally in our development because they comprise the smallest $\sigma$-algebra that contains sets of the form $\{f>t\}$ where $f \in C_{c}(X)$ and $t \in \mathbb{R}$. These are precisely the sets that occur in Theorem 9.2 when $L$ is taken as $C_{c}(X)$. In view of Exercise 2, Section 9.2, note that the outer measure obtained in the preceding Corollary is Borel. In general, it is possible to have access to Borel outer measures rather than merely Baire measures, as seen in the following result.
9.12. Theorem. Assume the hypotheses and notation of the Riesz Representation Theorem. Then there is a signed Borel outer measure $\bar{\mu}$ on $X$ with the property that

$$
\begin{equation*}
\int f d \bar{\mu}=T(f)=\int f d \mu \tag{9.12}
\end{equation*}
$$

for all $f \in C_{c}(X)$.

Proof. Assuming first that $T$ satisfies the conditions of Theorem 9.2, we will prove there is a Borel outer measure $\bar{\mu}$ such that (9.12) holds for every $f \in L^{+}$. Then referring to the proof of the Riesz Representation Theorem, it follows that there is a signed outer measure $\bar{\mu}$ that satisfies (9.12).

We define $\bar{\mu}$ in the following way. For every open set $U$ let

$$
\begin{equation*}
\alpha(U):=\sup \left\{T(f): f \in L^{+}, f \leq 1 \text { and } \operatorname{spt} f \subset U\right\} \tag{9.13}
\end{equation*}
$$

and for every $A \subset X$, define

$$
\bar{\mu}(A):=\inf \{\alpha(U): A \subset U, U \text { open }\}
$$

Observe that $\bar{\mu}$ and $\alpha$ agree on all open sets.
First, we verify that $\bar{\mu}$ is countably subadditive. Let $\left\{U_{i}\right\}$ be a sequence of open sets and let

$$
V=\bigcup_{i=1}^{\infty} U_{i} .
$$

If $f \in L^{+}, f \leq 1$, and $\operatorname{spt} f \subset V$, the compactness of $\operatorname{spt} f$ implies that

$$
\operatorname{spt} f \subset \bigcup_{i=1}^{N} U_{i}
$$

for some positive integer $N$. Now appeal to Lemma 11.20 to obtain functions $g_{1}, g_{2}, \ldots, g_{N} \in L^{+}$with $g_{i} \leq 1$, spt $g_{i} \subset U_{i}$, and

$$
\sum_{i=1}^{N} g_{i}(x)=1 \quad \text { whenever } \quad x \in \operatorname{spt} f
$$

Thus,

$$
f(x)=\sum_{i=1}^{N} g_{i}(x) f(x)
$$

and therefore

$$
T(f)=\sum_{i=1}^{N} T\left(g_{i} f\right) \leq \sum_{i=1}^{N} \alpha\left(U_{i}\right)<\sum_{i=1}^{\infty} \alpha\left(U_{i}\right)
$$

This implies

$$
\alpha\left(\bigcup_{i=1}^{\infty} U_{i}\right) \leq \sum_{i=1}^{\infty} \alpha\left(U_{i}\right) .
$$

From this and the definition of $\bar{\mu}$, it follows easily that $\bar{\mu}$ is countably subadditive on all sets.

Next, to show that $\bar{\mu}$ is a Borel outer measure, it suffices to show that each open set $U$ is $\bar{\mu}$ is measurable. For this, let $\varepsilon>0$ and let $A \subset X$ be an arbitrary set with $\bar{\mu}(A)<\infty$. Choose an open set $V \supset A$ such that $\alpha(V)<\bar{\mu}(A)+\varepsilon$ and a function $f \in L^{+}$with $f \leq 1$, spt $f \subset V \cap U$, and $T(f)>\alpha(V \cap U)-\varepsilon$. Also,
choose $g \in L^{+}$with $g \leq 1$ and $\operatorname{spt} g \subset V-\operatorname{spt} f$ so that $T(g)>\alpha(V-\operatorname{spt} f)-\varepsilon$. Then, since $f+g \leq 1$, we obtain

$$
\begin{aligned}
\bar{\mu}(A)+\varepsilon & \geq \alpha(V) \geq T(f+g)=T(f)+T(g) \\
& \geq \alpha(V \cap U)+\alpha(V-\operatorname{spt} f) \\
& \geq \bar{\mu}(A \cap U)+\bar{\mu}(A-U)-2 \varepsilon
\end{aligned}
$$

This shows that $U$ is $\bar{\mu}$-measurable since $\varepsilon$ is arbitrary.
Finally, to establish (9.12), by Theorem 6.60 it suffices to show that

$$
\bar{\mu}(\{f>s\})=\mu(\{f>s\})
$$

for all $f \in L^{+}$and all $s \in \mathbb{R}$. It follows from definitions that

$$
\bar{\mu}(U)=\alpha(U) \leq \mu(U)
$$

whenever $U$ is an open set. In particular, we have $\bar{\mu}\left(U_{s}\right) \leq \mu\left(U_{s}\right)$ where $U_{s}$ : = $\{f>s\}$. To prove the opposite inequality, note that if $f \in L^{+}$and $0<s<t$, then

$$
\alpha\left(U_{s}\right) \geq \frac{T[f \wedge(t+h)-f \wedge t]}{h} \text { for } \quad h>0
$$

because each

$$
f_{h}:=\frac{f \wedge(t+h)-f \wedge t}{h}= \begin{cases}1 & \text { if } f \geq t+h \\ \frac{f-t}{h} & \text { if } t<f<t+h \\ 0 & \text { if } f \leq t\end{cases}
$$

is a competitor in (9.13). Furthermore, since $f_{h_{1}} \geq f_{h_{2}}$ when $h_{1} \leq h_{2}$ and $\lim _{h \rightarrow 0^{+}} f_{h}=\chi_{\{f>t\}}$, we may apply the Monotone Convergence Theorem to obtain

$$
\begin{aligned}
\mu(\{f>t\}) & =\lim _{h \rightarrow 0^{+}} \int_{X} f_{h} d \mu \\
& =\lim _{h \rightarrow 0^{+}} \frac{T[f \wedge(t+h)-f \wedge t]}{h} \\
& \leq \alpha\left(U_{s}\right)=\bar{\mu}\left(U_{s}\right)
\end{aligned}
$$

Since

$$
\mu\left(U_{s}\right)=\lim _{t \rightarrow s^{+}} \mu(\{f>t\})
$$

we have $\mu\left(U_{s}\right) \leq \bar{\mu}\left(U_{s}\right)$.
From Corollary 9.4 we see that $\bar{\mu}$ and $\mu$ agree on compact sets and therefore, $\bar{\mu}$ is finite on compact sets. Furthermore, it follows from Lemma 9.3 that for an arbitrary set $A \subset X$, there exists a Borel set $B \supset A$ such that $\bar{\mu}(B)=\bar{\mu}(A)$. Recall that an outer measure with these properties is called a Radon outer measure
(see Definition 4.14). Note that this definition is compatible with that of Radon measure given in Definition 4.47.

## Exercises for Section 9.2

1. Provide an alternate (and simpler) proof of Urysohn's Lemma (Lemma 9.7) in the case when $X$ is a locally compact metric space.
2. Suppose $X$ is a locally compact Hausdorff space.
(i) If $f \in C_{c}(X)$ is nonnegative, then $\{a \leq f \leq \infty\}$ is a compact $G_{\delta}$ set for all $a>0$.
(ii) If $K \subset X$ is a compact $G_{\delta}$ set, there exists $f \in C_{c}(X)$ with $0 \leq f \leq 1$ and $K=f^{-1}(1)$.
(iii) The Baire sets are generated by compact $G_{\delta}$ sets.
3. Prove that the Baire sets and Borel sets are the same in a locally compact Hausdorff space satisfying the second axiom of countability.
4. Consider $(X, \mathcal{M}, \mu)$ where $X$ is a compact Hausdorff space and $\mathcal{M}$ is the family of Borel sets. If $\left\{\nu_{i}\right\}$ is a sequence of Radon measures with $\nu_{i}(X) \leq M$ for some positive $M$, then Alaoglu's Theorem (Theorem 8.40) implies there is a subsequence $\nu_{i_{j}}$ that converges weak* to some Radon measure $\nu$. Suppose $\left\{f_{i}\right\}$ is a sequence of functions in $L^{1}(X, \mathcal{M}, \mu)$. What can be concluded if $\left\|f_{i}\right\|_{1 ; \mu} \leq M$ ?
5. With the same notation as in the previous problem, suppose that $\nu_{i}$ converges weak ${ }^{*}$ to $\nu$. Prove the following:
(i) $\lim \sup _{i \rightarrow \infty} \nu_{i}(F) \leq \nu(F)$ whenever $F$ is a closed set.
(ii) $\liminf _{i \rightarrow \infty} \nu_{i}(U) \geq \nu(U)$ whenever $U$ is an open set.
(iii) $\lim _{i \rightarrow \infty} \nu_{i}(E)=\nu(E)$ whenever $E$ is a measurable set with $\nu(\partial E)=0$.
6. Assume that $X$ is a locally compact Hausdorff space that can be written as the countable union of compact sets. Then, under the hypotheses and notation of the Riesz Representation Theorem, prove that there is a (nonnegative) Baire outer measure $\mu$ and a $\mu$-measurable function $g$ such that
(i) $|g(x)|=1$ for $\mu$-a.e. $x$, and
(ii) $T(f)=\int_{X} f g d \mu$ for all $f \in C_{c}(X)$.

## CHAPTER 10

## Distributions

### 10.1. The Space $\mathcal{D}$

In the previous chapter, we saw how a bounded linear functional on the space $C_{c}\left(\mathbb{R}^{n}\right)$ can be identified with a measure. In this chapter, we will pursue this idea further by considering linear functionals on a smaller space, thus producing objects called distributions that are more general than measures. Distributions are of fundamental importance in many areas such as partial differential equations, the calculus of variations, and harmonic analysis. Their importance was formally acknowledged by the mathematics community when Laurent Schwartz, who initiated and developed the theory of distributions, was awarded the Fields Medal at the 1950 International Congress. In the next chapter some applications of distributions will be given. In particular, it will be shown how distributions are used to obtain a solution to a fundamental problem in partial differential equations, namely, the Dirichlet Problem. We begin by introducing the space of functions on which distributions are defined.

Let $\Omega \subset \mathbb{R}^{n}$ be an open set. We begin by investigating a space that is much smaller than $C_{c}(\Omega)$, namely, the space $\mathcal{D}(\Omega)$ of all infinitely differentiable functions whose supports are contained in $\Omega$. We let $C^{k}(\Omega), 1 \leq k \leq \infty$, denote the class of functions defined on $\Omega$ whose partial derivatives of all orders up to and including $k$ are continuous. Also, we denote by $C_{c}^{k}(\Omega)$ those functions in $C^{k}(\Omega)$ whose supports are contained in $\Omega$.

It is not immediately obvious that such functions exist, so we begin by analyzing the following function defined on $\mathbb{R}$ :

$$
f(x)= \begin{cases}e^{-1 / x} & x>0 \\ 0 & x \leq 0\end{cases}
$$

Observe that $f$ is $C^{\infty}$ on $\mathbb{R}-\{0\}$. It remains to show that all derivatives exist and are continuous at $x=0$. Now,

$$
f^{\prime}(x)= \begin{cases}\frac{1}{x^{2}} e^{-1 / x} & x>0 \\ 0 & x<0\end{cases}
$$

and therefore

$$
\begin{equation*}
\lim _{x \rightarrow 0} f^{\prime}(x)=0 \tag{10.1}
\end{equation*}
$$

Also,

$$
\begin{equation*}
\lim _{h \rightarrow 0^{+}} \frac{f(h)-f(0)}{h}=0 \tag{10.2}
\end{equation*}
$$

Note that (10.2) implies that $f^{\prime}(0)$ exists and $f^{\prime}(0)=0$. Moreover, (10.1) gives that $f^{\prime}$ is continuous at $x=0$.

A similar argument establishes the same conclusion for all higher derivatives of $f$. Indeed, a direct calculation shows that the $k^{\text {th }}$ derivatives of $f$ for $x \neq 0$ are of the form $P_{2 k}\left(\frac{1}{x}\right) f(x)$, where $P$ is a polynomial of degree $2 k$. Thus

$$
\begin{align*}
\lim _{x \rightarrow 0^{+}} f^{(k)}(x) & =\lim _{x \rightarrow 0^{+}} P\left(\frac{1}{x}\right) e^{\frac{-1}{x}} .  \tag{10.3}\\
& =\lim _{x \rightarrow 0^{+}} \frac{P\left(\frac{1}{x}\right)}{e^{\frac{1}{x}}} \\
& =\lim _{t \rightarrow \infty} \frac{P(t)}{e^{t}}=0
\end{align*}
$$

where this limit is computed by using L'Hôpital's rule repeatedly. Finally,

$$
\begin{align*}
\lim _{h \rightarrow 0^{+}} \frac{f^{(k-1)}(h)-f^{(k-1)}(0)}{h} & =\lim _{h \rightarrow 0^{+}} \frac{f^{(k-1)}(h)}{h}  \tag{10.4}\\
& =\lim _{h \rightarrow 0^{+}} \frac{P_{2(k-1)}\left(\frac{1}{h}\right) e^{\frac{-1}{h}}}{h} \\
& =\lim _{h \rightarrow 0^{+}} Q\left(\frac{1}{h}\right) e^{\frac{-1}{h}} \\
& =\lim _{h \rightarrow 0^{+}} \frac{Q(t)}{e^{t}}=0 .
\end{align*}
$$

where $Q$ is a polynomial of degree $2 k-1$. From (10.4) we conclude $f^{(k)}$ exists at $x=0$ and $f^{(k)}(0)=0$. Moreover, (10.3) shows that $f^{(k)}$ is continuous at $x=0$.

We can now construct a $C^{\infty}$ function with compact support in $\mathbb{R}^{n}$. For this, let

$$
F(x)=f\left(1-|x|^{2}\right), x \in \mathbb{R}^{n} .
$$

With $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, observe that $1-|x|^{2}$ is the polynomial $1-\left(x_{1}^{2}+x_{2}^{2}+\cdots+\right.$ $x_{n}^{2}$ ) and therefore, that $F$ is an infinitely differentiable function of $x$. Moreover, $F$ is nonnegative and is zero for $|x| \geq 1$.

It is traditional in many parts of analysis to denote $C^{\infty}$ functions with compact support by $\varphi$. We will adopt this convention. Now for some notation. The partial derivative operators are denoted by $D_{i}=\partial / \partial x_{i}$ for $1 \leq i \leq n$. Thus,

$$
\begin{equation*}
D_{i} \varphi=\frac{\partial \varphi}{\partial x_{i}} . \tag{10.5}
\end{equation*}
$$

If $\alpha=\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right)$ is an $n$-tuple of nonnegative integers, $\alpha$ is called a multiindex and the length of $\alpha$ is defined as

$$
|\alpha|=\sum_{i=1}^{n} \alpha_{i}
$$

Using this notation, higher order derivatives are denoted by

$$
D^{\alpha} \varphi=\frac{\partial^{|\alpha|} \varphi}{\partial x_{1}^{\alpha_{1}} \ldots \partial x_{n}^{\alpha_{n}}}
$$

For example,

$$
\frac{\partial^{2} \varphi}{\partial x \partial y}=D^{(1,1)} \varphi \quad \text { and } \quad \frac{\partial^{3} \varphi}{\partial^{2} x \partial y}=D^{(2,1)} \varphi
$$

Also, we let $\nabla \varphi(x)$ denote the gradient of $\varphi$ at $x$, that is,

$$
\begin{align*}
\nabla \varphi(x) & =\left(\frac{\partial \varphi}{\partial x_{1}}(x), \frac{\partial \varphi}{\partial x_{2}}(x), \ldots, \frac{\partial \varphi}{\partial x_{n}}(x)\right)  \tag{10.6}\\
& =\left(D_{1} \varphi(x), D_{2} \varphi(x), \ldots, D_{n} \varphi(x)\right)
\end{align*}
$$

We have shown that there exist $C^{\infty}$ functions $\varphi$ with the property that $\varphi(x)>0$ for $x \in B(0,1)$ and that $\varphi(x)=0$ whenever $|x| \geq 1$. By multiplying $\varphi$ by a suitable constant, we can assume

$$
\begin{equation*}
\int_{B(0,1)} \varphi(x) d \lambda(x)=1 \tag{10.7}
\end{equation*}
$$

By employing an appropriate scaling of $\varphi$, we can duplicate these properties on the ball $B(0, \varepsilon)$ for any $\varepsilon>0$. For this purpose, let

$$
\varphi_{\varepsilon}(x)=\varepsilon^{-n} \varphi\left(\frac{x}{\varepsilon}\right)
$$

Then $\varphi_{\varepsilon}$ has the same properties on $B(0, \varepsilon)$ as does $\varphi$ on $B(0,1)$. Furthermore,

$$
\begin{equation*}
\int_{B(0, \varepsilon)} \varphi_{\varepsilon}(x) d \lambda(x)=1 \tag{10.8}
\end{equation*}
$$

Thus, we have shown that $\mathcal{D}(\Omega)$ is nonempty; we will often call functions in $\mathcal{D}(\Omega)$ test functions.

We now show how the functions $\varphi_{\varepsilon}$ can be used to generate more $C^{\infty}$ functions with compact support. Given a function $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$, recall the definition of convolution first introduce in Section 6.11:

$$
\begin{equation*}
f * \varphi_{\varepsilon}(x)=\int_{\mathbb{R}^{n}} f(x-y) \varphi_{\varepsilon}(y) d \lambda(y) \tag{10.9}
\end{equation*}
$$

for all $x \in \mathbb{R}^{n}$. We will use the notation $f_{\varepsilon}:=f * \varphi_{\varepsilon}$. The function $f_{\varepsilon}$ is called the mollifier of $f$.
10.1. Theorem.
(i) If $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$, then for every $\varepsilon>0, f_{\varepsilon} \in C^{\infty}\left(\mathbb{R}^{n}\right)$.
(ii)

$$
\lim _{\varepsilon \rightarrow 0} f_{\varepsilon}(x)=f(x)
$$

whenever $x$ is a Lebesgue point for $f$. In case $f$ is continuous, then $f_{\varepsilon}$ converges to $f$ uniformly on compact subsets of $\mathbb{R}^{n}$.
(iii) If $f \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<\infty$, then $f_{\varepsilon} \in L^{p}\left(\mathbb{R}^{n}\right),\left\|f_{\varepsilon}\right\|_{p} \leq\|f\|_{p}$, and $\lim _{\varepsilon \rightarrow 0} \| f_{\varepsilon}-$ $f \|_{p}=0$.

Proof. Recall that $D_{i} f$ denotes the partial derivative of $f$ with respect to the $i^{\text {th }}$ variable. The proof that $f_{\varepsilon}$ is $C^{\infty}$ will be established if we can show that

$$
\begin{equation*}
D_{i}\left(\varphi_{\varepsilon} * f\right)=\left(D_{i} \varphi_{\varepsilon}\right) * f \tag{10.10}
\end{equation*}
$$

for each $i=1,2, \ldots, n$. To see this, assume (10.10) for the moment. The right side of the equation is a continuous function (Exercise 2, Section 10.1), thus showing that $f_{\varepsilon}$ is $C^{1}$. Then, if $\alpha$ denotes the $n$-tuple with $|\alpha|=2$ and with 1 in the $i^{\text {th }}$ and $j^{\text {th }}$ positions, it follows that

$$
\begin{aligned}
D^{\alpha}\left(\varphi_{\varepsilon} * f\right) & =D_{i}\left[D_{j}\left(\varphi_{\varepsilon} * f\right)\right] \\
& =D_{i}\left[\left(D_{j} \varphi_{\varepsilon}\right) * f\right] \\
& =\left(D^{\alpha} \varphi_{\varepsilon}\right) * f,
\end{aligned}
$$

which proves that $f_{\varepsilon} \in C^{2}$ since again the right side of the equation is continuous. Proceeding this way by induction, it follows that $f \in C^{\infty}$.

We now turn to the proof of (10.10). Let $e_{1}, \ldots, e_{n}$ be the standard basis of $\mathbb{R}^{n}$ and consider the partial derivative with respect to the $i^{\text {th }}$ variable. For every real number $h$, we have

$$
f_{\varepsilon}\left(x+h e_{i}\right)-f_{\varepsilon}(x)=\int_{\mathbb{R}^{n}}\left[\varphi_{\varepsilon}\left(x-z+h e_{i}\right)-\varphi_{\varepsilon}(x-z)\right] f(z) d \lambda(z)
$$

Let $\alpha(t)$ denote the integrand:

$$
\alpha(t)=\varphi_{\varepsilon}\left(x-z+t e_{i}\right) f(z)
$$

Then $\alpha$ is a $C^{\infty}$ function of $t$. The chain rule implies

$$
\begin{aligned}
\alpha^{\prime}(t) & =\nabla \varphi_{\varepsilon}\left(x-z+t e_{i}\right) \cdot e_{i} f(z) \\
& =D_{i} \varphi_{\varepsilon}\left(x-z+t e_{i}\right) f(z)
\end{aligned}
$$

Since

$$
\alpha(h)-\alpha(0)=\int_{0}^{h} \alpha^{\prime}(t) d t
$$

we have

$$
\begin{aligned}
& f_{\varepsilon}\left(x+h e_{i}\right)-f_{\varepsilon}(x) \\
& =\int_{\mathbb{R}^{n}}\left[\varphi_{\varepsilon}\left(x-z+h e_{i}\right)-\varphi_{\varepsilon}(x-z)\right] f(z) d \lambda(z) \\
& =\int_{\mathbb{R}^{n}} \alpha(h)-\alpha(0) d \lambda(z) \\
& =\int_{\mathbb{R}^{n}} \int_{0}^{h} D_{i} \varphi_{\varepsilon}\left(x-z+t e_{i}\right) f(z) d t d \lambda(z) \\
& =\int_{0}^{h} \int_{\mathbb{R}^{n}} D_{i} \varphi_{\varepsilon}\left(x-z+t e_{i}\right) f(z) d \lambda(z) d t .
\end{aligned}
$$

It follows from Lebesgue's Dominated Convergence Theorem that the inner integral is a continuous function of $t$. Now divide both sides of the equation by $h$ and take the limit as $h \rightarrow 0$ to obtain

$$
D_{i} f_{\varepsilon}(x)=\int_{\mathbb{R}^{n}} D_{i} \varphi_{\varepsilon}(x-z) f(z) d \lambda(z)=\left(D_{i} \varphi_{\varepsilon}\right) * f(x)
$$

which establishes (10.10) and therefore (i).
In case (ii), observe that

$$
\begin{align*}
& \left|f_{\varepsilon}(x)-f(x)\right| \\
& \leq \int_{\mathbb{R}^{n}} \varphi_{\varepsilon}(x-y)|f(y)-f(x)| d \lambda(y)  \tag{10.11}\\
& \leq \max _{\mathbb{R}^{n}} \varphi \varepsilon^{-n} \int_{B(x, \varepsilon)}|f(x)-f(y)| d \lambda(y) \rightarrow 0
\end{align*}
$$

as $\varepsilon \rightarrow 0$ whenever $x$ is a Lebesgue point for $f$. Now consider the case where $f$ is continuous and $K \subset \mathbb{R}^{n}$ is a compact. Then $f$ is uniformly continuous on $K$ and also on the closure of the open set $U=\{x: \operatorname{dist}(x, K)<1\}$. For each $\eta>0$, there exists $0<\varepsilon<1$ such that $|f(x)-f(y)|<\eta$ whenever $|x-y|<\varepsilon$ and whenever $x, y \in U$; in particular when $x \in K$. Consequently, it follows from (10.11) that whenever $x \in K$,

$$
\left|f_{\varepsilon}(x)-f(x)\right| \leq M \eta
$$

where $M$ is the product of $\max _{\mathbb{R}^{n}} \varphi$ and the Lebesgue measure of the unit ball. Since $\eta$ is arbitrary, this shows that $f_{\varepsilon}$ converges uniformly to $f$ on $K$.

The first part of (iii) follows from Theorem 6.57 since $\left\|\varphi_{\varepsilon}\right\|_{1}=1$.
Finally, addressing the second part of (iii), for each $\eta>0$, select a continuous function $g$ with compact support on $\mathbb{R}^{n}$ such that

$$
\begin{equation*}
\|f-g\|_{p}<\eta \tag{10.12}
\end{equation*}
$$

(see Exercise 12, Section 6.5). Because $g$ has compact support, it follows from (ii) that $\left\|g-g_{\varepsilon}\right\|_{p}<\eta$ for all $\varepsilon$ sufficiently small. Now apply Theorem 10.1 (iii) and (10.12) to the difference $g-f$ to obtain

$$
\left\|f-f_{\varepsilon}\right\|_{p} \leq\|f-g\|_{p}+\left\|g-g_{\varepsilon}\right\|_{p}+\left\|g_{\varepsilon}-f_{\varepsilon}\right\|_{p} \leq 3 \eta
$$

This shows that $\left\|f-f_{\varepsilon}\right\|_{p} \rightarrow 0$ as $\varepsilon \rightarrow 0$.

## Exercises for 10.1

1. Let $K$ be a compact subset of an open set $\Omega$. Prove that there exists $f \in C_{c}^{\infty}(\Omega)$ such that $f \equiv 1$ on $K$.
2. Suppose $f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ and $\varphi$ a continuous function with compact support. Prove that $\varphi * f$ is continuous.
3. If $f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ and

$$
\int_{\mathbb{R}^{n}} f \varphi d x=0
$$

for every $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, show that $f=0$ almost everywhere. Hint: use Theorem 10.1.

### 10.2. Basic Properties of Distributions

10.2. Definition. Let $\Omega \subset \mathbb{R}^{n}$ be an open set. A linear functional $T$ on $\mathcal{D}(\Omega)$ is a distribution if and only if for every compact set $K \subset \Omega$, there exist constants $C$ and $N$ such that

$$
|T(\varphi)| \leq C(K) \sup _{x \in K} \sum_{|\alpha| \leq N(K)}\left|D^{\alpha} \varphi(x)\right|
$$

for all test functions $\varphi \in \mathcal{D}(\Omega)$ with support in $K$. If the integer $N$ can be chosen independent of the compact set $K$, and $N$ is the smallest possible choice, the distribution is said to be of order $N$. We use the notation

$$
\|\varphi\|_{K ; N}:=\sup _{x \in K} \sum_{|\alpha| \leq N}\left|D^{\alpha} \varphi(x)\right| .
$$

Thus, in particular, $\|\varphi\|_{K ; 0}$ denotes the sup norm of $\varphi$ on $K$.
Here are some examples of distributions. First, suppose $\mu$ is a signed Radon measure on $\Omega$. Define the corresponding distribution by

$$
T(\varphi)=\int_{\Omega} \varphi(x) d \mu(x)
$$

for every test function $\varphi$. Note that the integral is finite since $\mu$ is finite on compact sets, by definition. If $K \subset \Omega$ is compact, and $C(K)=|\mu|(K)$ (recall the notation
in 6.39), then

$$
\begin{aligned}
|T(\varphi)| & \leq C(K)\|\varphi\|_{L^{\infty}} \\
& =C(K)\|\varphi\|_{K ; 0}
\end{aligned}
$$

for all test functions $\varphi$ with support in $K$. Thus, the distribution $T$ corresponding to the measure $\mu$ is of order 0 . In the context of distribution theory, a Radon measure will be identified as a distribution in this way. In particular, consider the Dirac measure $\delta$ whose total mass is concentrated at the origin:

$$
\delta(E)= \begin{cases}1 & 0 \in E \\ 0 & 0 \notin E\end{cases}
$$

The distribution identified with this measure is defined by

$$
\begin{equation*}
T(\varphi)=\varphi(0) \tag{10.13}
\end{equation*}
$$

for every test function $\varphi$.
10.3. Definition. Let $f \in L_{\mathrm{loc}}^{1}(\Omega)$. The distribution corresponding to $f$ is defined as

$$
T(\varphi)=\int_{\Omega} \varphi(x) f(x) d \lambda(x)
$$

Thus, a locally integrable function can be considered as an absolutely continuous measure and is therefore identified with a distribution of order 0 . The distribution corresponding to $f$ is sometimes denoted as $T_{f}$.

We have just seen that a Radon measure is a distribution of order 0. The following result shows that we can actually identify measures and distributions of order 0 .
10.4. Theorem. A distribution $T$ is a Radon measure if and only if $T$ is of order 0 .

Proof. Assume that $T$ is a distribution of order 0 . We will show that $T$ can be extended in a unique way to a linear functional $T^{*}$ on the space $C_{c}(\Omega)$ so that

$$
\begin{equation*}
\left|T^{*}(\varphi)\right| \leq C(K)\|\varphi\|_{K ; 0} \tag{10.14}
\end{equation*}
$$

for each compact set $K$ and $\varphi$ supported on $K$. Then, by appealing to the Riesz Representation Theorem, it will follow that there is a (signed) Radon measure $\mu$ such that

$$
T^{*}(\varphi)=\int_{\mathbb{R}^{n}} \varphi d \mu
$$

for each $\varphi \in C_{c}(\Omega)$. In particular, we will have

$$
T(\varphi)=T^{*}(\varphi)=\int_{\mathbb{R}^{n}} \varphi d \mu
$$

whenever $\varphi \in \mathcal{D}(\Omega)$, thus establishing that $\mu$ is the measure identified with $T$.
In order to prove (10.14), select a continuous function $\varphi$ with support in a compact set $K$. By using mollifiers and Theorem 10.1, it follows that there is a sequence of test functions $\left\{\varphi_{i}\right\} \in \mathcal{D}(\Omega)$ whose supports are contained in a fixed compact neighborhood of $\operatorname{spt} \varphi$ such that

$$
\left\|\varphi_{i}-\varphi\right\|_{K ; 0} \rightarrow 0 \quad \text { as } \quad i \rightarrow \infty .
$$

Now define

$$
T^{*}(\varphi)=\lim _{i \rightarrow \infty} T\left(\varphi_{i}\right) .
$$

The limit exists because when $i, j \rightarrow \infty$, we have

$$
\left|T\left(\varphi_{i}\right)-T\left(\varphi_{j}\right)\right|=\left|T\left(\varphi_{i}-\varphi_{j}\right)\right| \leq C\left\|\varphi_{i}-\varphi_{j}\right\|_{K ; 0} \rightarrow 0 .
$$

Similar reasoning shows that the limit is independent of the sequence chosen. Furthermore, since $T$ is of order 0 , it follows that

$$
\left|T^{*}(\varphi)\right| \leq C\|\varphi\|_{K ; 0},
$$

which establishes (10.14).
We conclude this section with a simple but very useful condition that ensures that a distribution is a measure.
10.5. Definition. A distribution on $\Omega$ is positive if $T(\varphi) \geq 0$ for all test functions on $\Omega$ satisfying $\varphi \geq 0$.
10.6. Theorem. A distribution $T$ on $\Omega$ is positive if and only if $T$ is a positive measure.

Proof. From previous discussions, we know that a Radon measure is a distribution is of order 0 , so we need only consider the case when $T$ is a positive distribution. Let $K \subset \Omega$ be a compact set. From Exercise 1, Section 10.1, there exists a function $\alpha \in C_{c}^{\infty}(\Omega)$ that equals 1 on a neighborhood of $K$. Now select a test function $\varphi$ whose support is contained in $K$. Then we have

$$
-\|\varphi\|_{K ; 0} \alpha(x) \leq \varphi(x) \leq\|\varphi\|_{K ; 0} \alpha(x) \quad \text { for all } x,
$$

and therefore

$$
-\|\varphi\|_{K ; 0} T(\alpha) \leq T(\varphi) \leq\|\varphi\|_{K ; 0} T(\alpha) .
$$

Thus, $|T(\varphi)| \leq|T(\alpha)|\|\varphi\|_{K ; 0}$, which shows that $T$ is of order 0 and therefore a measure $\mu$. The measure $\mu$ is clearly non-negative.

## Exercises for Section 10.2

1. Use the Hahn-Banach theorem to provide an alternative proof of (371.1).
2. Show that the principal value integral

$$
\text { p.v. } \int \frac{\phi(x)}{x} d x=\lim _{\varepsilon \rightarrow 0^{+}}\left(\int_{-\infty}^{-\varepsilon} \frac{\phi(x)}{x} d x+\int_{\varepsilon}^{\infty} \frac{\phi(x)}{x} d x\right)
$$

exists for all $\phi \in \mathcal{D}(\mathbb{R})$, and is a distribution. What is its order?
3. Let $T_{j}, j=1,2, \ldots$ be a sequence of distributions. The sequence $T_{j}$ is said to converge to the distribution $T$ if

$$
\lim T_{j}(\varphi)=T(\varphi), \text { for all } \varphi \in \mathcal{D}(\Omega)
$$

Let $f_{\varepsilon} \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$ be a function which depends on a parameter $\varepsilon \in(0,1)$, and is such that
(a) $\operatorname{supp} f_{\varepsilon} \subset\{|x| \leq \varepsilon\}$
(b) $\int_{\mathbb{R}^{n}} f_{\varepsilon}(x) d x=1$
(c) $\int_{\mathbb{R}^{n}}\left|f_{\varepsilon}(x)\right| d x \leq \mu<\infty, 0<\varepsilon<1$.

Show that $f_{\varepsilon} \rightarrow \delta$ as $\varepsilon \rightarrow 0$. Show also that, if $f_{\varepsilon}$ satisfies (a) and $f_{\varepsilon} \rightarrow \delta$ as $\varepsilon \rightarrow 0$, then (b) holds.

### 10.3. Differentiation of Distributions

One of the primary reasons why distributions were created was to provide a notion of differentiability for functions that are not differentiable in the classical sense. In this section, we define the derivative of a distribution and investigate some of its properties.

In order to motivate the definition of a distribution, consider the following simple case. Let $\Omega$ denote the open interval $(0,1)$ in $\mathbb{R}$, and let $f$ denote an absolutely continuous function on $(0,1)$. If $\varphi$ is a test function in $\Omega$, we can integrate by parts (Exercise 7, Section 7.5) to obtain

$$
\begin{equation*}
\int_{0}^{1} f^{\prime}(x) \varphi(x) d \lambda(x)=-\int_{0}^{1} f(x) \varphi^{\prime}(x) d \lambda(x) \tag{10.15}
\end{equation*}
$$

since by definition, $\varphi(1)=\varphi(0)=0$. If we consider $f$ as a distribution $T$, we have

$$
T(\varphi)=\int_{0}^{1} f(x) \varphi(x) d \lambda(x)
$$

for every test function $\varphi$ in $(0,1)$. Now define a distribution $S$ by

$$
S(\varphi)=\int_{0}^{1} f(x) \varphi^{\prime}(x) d \lambda(x)
$$

and observe that it is a distribution of order 1. From (10.15) we see that $S$ can be identified with $-f^{\prime}$. From this it is clear that the derivative $T^{\prime}$ should be defined
as

$$
T^{\prime}(\varphi)=-T\left(\varphi^{\prime}\right)
$$

for every test function $\varphi$. More generally, we have the following definition.
10.7. Definition. Let $T$ be a distribution of order $N$ defined on an open set $\Omega \subset \mathbb{R}^{n}$. The partial derivative of $T$ with respect to the $i^{\text {th }}$ coordinate direction is defined by

$$
\frac{\partial T}{\partial x_{i}}(\varphi)=-T\left(\frac{\partial \varphi}{\partial x_{i}}\right) .
$$

Observe that since the derivative of a test function is again a test function, the differentiated distribution is a linear functional on $\mathcal{D}(\Omega)$. It is in fact a distribution since

$$
\left|T\left(\frac{\partial \varphi}{\partial x_{i}}\right)\right| \leq C\|\varphi\|_{N+1}
$$

is valid whenever $\varphi$ is a test function supported by a compact set $K$ on which

$$
|T(\varphi)| \leq C\|\varphi\|_{N}
$$

Let $\alpha$ be any multi-index. More generally, the $\alpha^{t h}$ derivative of the distribution $T$ is another distribution defined by:

$$
D^{\alpha} T(\varphi)=(-1)^{|\alpha|} T\left(D^{\alpha} \varphi\right)
$$

Recall that a function $f \in L_{l o c}^{1}(\Omega)$ is associated with the distribution $T_{f}$ (see Definition 10.3). Thus, if $f, g \in L_{l o c}^{1}(\Omega)$ we say that $D^{\alpha} f=g_{\alpha}$ in the sense of distributions if

$$
\int_{\Omega} f D^{\alpha} \varphi d \lambda=(-1)^{|\alpha|} \int_{\Omega} \varphi g_{\alpha} d \lambda, \text { for all } \varphi \in \mathcal{D}(\Omega)
$$

Let us consider some examples. The first one has already been discussed above, but is repeated for emphasis.

1. Let $\Omega=(a, b)$, and suppose $f$ is an absolutely continuous function defined on $[a, b]$. If $T$ is the distribution corresponding to $f$, we have

$$
T(\varphi)=\int_{a}^{b} \varphi f d \lambda
$$

for each test function $\varphi$. Since $f$ is absolutely continuous and $\varphi$ has compact support in $[a, b]$ (so that $\varphi(b)=\varphi(a)=0$ ), we may employ integration by parts to conclude

$$
T^{\prime}(\varphi)=-T\left(\varphi^{\prime}\right)=-\int_{a}^{b} \varphi^{\prime} f d \lambda=\int_{a}^{b} \varphi f^{\prime} d \lambda
$$

Thus, the distribution $T^{\prime}$ is identified with the function $f^{\prime}$.
The next example shows how it is possible for a function to have a derivative in the sense of distributions, but not be differentiable in the classical sense.
2. Let $\Omega=\mathbb{R}$ and define

$$
f(x)= \begin{cases}1 & x>0 \\ 0 & x \leq 0\end{cases}
$$

With $T$ defined as the distribution corresponding to $f$, we obtain

$$
T^{\prime}(\varphi)=-T\left(\varphi^{\prime}\right)=-\int_{\mathbb{R}} \varphi^{\prime} f d \lambda=-\int_{0}^{\infty} \varphi^{\prime} d \lambda=\varphi(0)
$$

Thus, the derivative $T^{\prime}$ is equal to the Dirac measure, see (10.13).
3. We alter the function of Example 2 slightly:

$$
f(x)=|x| .
$$

Then

$$
T(\varphi)=\int_{0}^{\infty} x \varphi(x) d \lambda(x)-\int_{-\infty}^{0} x \varphi(x) d \lambda(x)
$$

so that, after integrating by parts, we obtain

$$
\begin{aligned}
T^{\prime}(\varphi) & =\int_{0}^{\infty} \varphi(x) d \lambda(x)-\int_{-\infty}^{0} \varphi(x) d \lambda(x) \\
& =\int_{\mathbb{R}} \varphi(x) g(x) d \lambda(x)
\end{aligned}
$$

where

$$
g(x):= \begin{cases}1 & x>0 \\ -1 & x \leq 0\end{cases}
$$

This shows that the derivative of $f$ is $g$ in the sense of distributions.
One would hope that the basic results of calculus carry over within the framework of distributions. The following is the first of many that do.
10.8. Theorem. If $T$ is a distribution in $\mathbb{R}$ with $T^{\prime}=0$, then $T$ is a constant. That is, $T$ is the distribution that corresponds to a constant function.

Proof. Observe that $\varphi=\psi^{\prime}$ where

$$
\psi(x):=\int_{-\infty}^{x} \varphi(t) d t
$$

Since $\varphi$ has compact support, it follows that $\psi$ has compact support precisely when

$$
\begin{equation*}
\int_{-\infty}^{\infty} \varphi(t) d t=0 \tag{10.16}
\end{equation*}
$$

Thus, $\varphi$ is the derivative of another test function if and only if (10.16) holds.
To prove the theorem, choose an arbitrary test function $\psi$ with the property

$$
\int \psi(x) d \lambda(x)=1
$$

Then any test function $\varphi$ can be written as

$$
\varphi(x)=[\varphi(x)-a \psi(x)]+a \psi(x)
$$

where

$$
a=\int \varphi(t) d \lambda(t)
$$

The function in brackets is a test function, call it $\alpha$, whose integral is 0 and is therefore the derivative of another test function. Since $T^{\prime}=0$, it follows that $T(\alpha)=0$ and therefore we obtain

$$
T(\varphi)=a T(\psi)=\int T(\psi) \varphi(t) d \lambda(t)
$$

which shows that $T$ corresponds to the constant $T(\psi)$.

This result along with Example 1 gives an interesting characterization of absolutely continuous functions.
10.9. Theorem. Suppose $f \in L_{\mathrm{loc}}^{1}(a, b)$. Then $f$ is equal almost everywhere to an absolutely continuous function on $[a, b]$ if and only if the derivative of the distribution corresponding to $f$ is a function.

Proof. In Example 1 we have already seen that if $f$ is absolutely continuous, then its derivative in the sense of distributions is again a function. Indeed, the function associated with the derivative of the distribution is $f^{\prime}$.

Now suppose that $T$ is the distribution associated with $f$ and that $T^{\prime}=g$. In order to show that $f$ is equal almost everywhere to an absolutely continuous function, let

$$
h(x)=\int_{a}^{x} g(t) d \lambda(t)
$$

Observe that $h$ is absolutely continuous and that $h^{\prime}=g$ almost everywhere. Let $S$ denote the distribution corresponding to $h$; that is,

$$
S(\varphi)=\int h \varphi
$$

for every test function $\varphi$. Then

$$
S^{\prime}(\varphi)=-\int h \varphi^{\prime}=\int h^{\prime} \varphi=\int g \varphi
$$

Thus, $S^{\prime}=g$. Since $T^{\prime}=g$ also, we have $T=S+k$ for some constant $k$ by the previous theorem. This implies that

$$
\begin{aligned}
\int f \varphi d \lambda=T(\varphi) & =S(\varphi)+\int k \varphi \\
& =\int h \varphi+k \varphi \\
& =\int(h+k) \varphi
\end{aligned}
$$

for every test function $\varphi$. This implies that $f=h+k$ almost everywhere (see Exercise 3, Section 10.1).

Since functions of bounded variation are closely related to absolutely continuous functions, it is natural to inquire whether they too can be characterized in terms of distributions.
10.10. Theorem. Suppose $f \in L_{\mathrm{loc}}^{1}(a, b)$. Then $f$ is equivalent to a function of bounded variation if and only if the derivative of the distribution corresponding to $f$ is a signed measure whose total variation is finite.

Proof. Suppose $f \in L_{\mathrm{loc}}^{1}(a, b)$ is a nondecreasing function and let $T$ be the distribution corresponding to $f$. Then, for every test function $\varphi$,

$$
T(\varphi)=\int f \varphi d \lambda
$$

and

$$
T^{\prime}(\varphi)=-T\left(\varphi^{\prime}\right)=-\int f \varphi^{\prime} d \lambda
$$

Since $f$ is nondecreasing, it generates a Lebesgue-Stieltjes measure $\lambda_{f}$. By the integration by parts formula for such measures (see Exercises 5 and 6 in Section 6.3 ), we have

$$
-\int f \varphi^{\prime} d \lambda=\int \varphi d \lambda_{f}
$$

which shows that $T^{\prime}$ corresponds to the measure $\lambda_{f}$. In case $f$ is of bounded variation on $(a, b)$, we can write

$$
f=f_{1}-f_{2}
$$

where $f_{1}$ and $f_{2}$ are nondecreasing functions. Then the distribution corresponding to $f, T_{f}$, can be written as $T_{f}=T_{f_{1}}-T_{f_{2}}$ and therefore

$$
T_{f}^{\prime}=T_{f_{1}}^{\prime}-T_{f_{2}}^{\prime}=\lambda_{f_{1}}-\lambda_{f_{2}}
$$

Consequently, the signed measure $\lambda_{f_{1}}-\lambda_{f_{2}}$ corresponds to the distributional derivative of $f$.

Conversely, suppose the derivative of $f$ is a signed measure $\mu$. By the Jordan Decomposition Theorem, we can write $\mu$ as the difference of two nonnegative measures, $\mu=\mu_{1}-\mu_{2}$. Let

$$
f_{i}(x)=\mu_{i}((-\infty, x]) \quad i=1,2
$$

From Theorem 4.33, we have that $\mu_{i}$ agrees with the Lebesgue-Stieltjes measure $\lambda_{f_{i}}$ on all Borel sets. Furthermore, utilizing the formula for integration by parts, we obtain

$$
T_{f_{i}}^{\prime}(\varphi)=-T_{f_{i}}\left(\varphi^{\prime}\right)=-\int f_{i} \varphi^{\prime} d \lambda=\int \varphi d \lambda_{f_{i}}=\int \varphi d \mu_{i}
$$

Thus, with $g:=f_{1}-f_{2}$, we have that $g$ is of bounded variation and that its distributional derivative is $\mu_{1}-\mu_{2}=\mu$. Since $f^{\prime}=\mu$, we conclude from Theorem 10.8 that $f-g$ is a constant, and therefore that $f$ is equivalent to a function of bounded variation.

## Exercises for Section 10.3

1. If a distribution $T$ defined on $\mathbb{R}$ has the property that $T^{\prime \prime}=0$, what can be said about $T$ ?
2. Show that if $T_{k}, k=1,2, \ldots$ is a sequence of distributions that converges to a distribution $T$ (see Exercise 3, Section 10.2), and $\alpha$ is a multi-index, then

$$
D^{\alpha} T_{k} \rightarrow D^{\alpha} T
$$

3. Let $\Omega=\mathbb{R}$ and define

$$
H(x)= \begin{cases}1 & x>0 \\ 0 & x \leq 0\end{cases}
$$

Let $T$ be the distribution corresponding to $H$, and $g \in C^{\infty}(\mathbb{R})$. Show that the function $x \rightarrow g(x) H(x)$ is locally integrable and so determines a distribution satisfying $D(g H)=g(0) \delta+H D g$.
4. Show that if $T$ is a distribution on $\mathcal{D}(\mathbb{R})$, and $x D T+T=0$, then $T=A\left(\frac{1}{x}\right)+B \delta$, where $A$ and $B$ are real numbers, and $\frac{1}{x}$ is the principal value distribution introduced in Exercise 2, Section 10.2.

### 10.4. Essential Variation

We have seen that a function of bounded variation defines a distribution whose derivative is a measure. Question: How is the variation of the function related to the variation of the measure? The notion of essential variation provides the answer.

A function $f$ of bounded variation gives rise to a distribution whose derivative is a measure. Furthermore, any other function $g$ agreeing with $f$ almost everywhere
defines the same distribution. Of course, $g$ need not be of bounded variation. This raises the question of what condition on $f$ is equivalent to $f^{\prime}$ being a measure. We will see that the needed ingredient is a concept of variation that remains unchanged when the function is altered on a set of measure zero.
10.11. Definition. The essential variation of a function $f$ defined on $(a, b)$ is

$$
\operatorname{ess} V_{a}^{b} f=\sup \left\{\sum_{i=1}^{k}\left|f\left(t_{i+1}\right)-f\left(t_{i}\right)\right|\right\}
$$

where the supremum is taken over all finite partitions $a<t_{1}<\cdots<t_{k+1}<b$ such that each $t_{i}$ is a point of approximate continuity of $f$.

Recall that the variation of a function is defined similarly, but without the restriction that $f$ be approximately continuous at each $t_{i}$. Clearly, if $f=g$ a.e. on $(a, b)$, then

$$
\operatorname{ess} V_{a}^{b} f=\operatorname{ess} V_{a}^{b} g
$$

A signed Radon measure $\mu$ on $(a, b)$ defines a linear functional on the space of continuous functions with compact support in $(a, b)$ by

$$
T_{\mu}(\varphi)=\int_{a}^{b} \varphi d \mu
$$

Recall that the total variation of $\mu$ on $(a, b)$ is defined (see definition 6.39) as the norm of $T_{\mu}$; that is,

$$
\begin{equation*}
\|\mu\|=\sup \left\{\int_{a}^{b} \varphi d \mu: \varphi \in C_{c}(a, b),|\varphi| \leq 1\right\} \tag{10.17}
\end{equation*}
$$

Notice that the supremum could just as well be taken over $C^{\infty}$ functions with compact support since they are dense in $C_{c}(a, b)$ in the topology of uniform convergence.
10.12. Theorem. Suppose $f \in L^{1}(a, b)$. Then $f^{\prime}$ (in the sense of distributions) is a measure with finite total variation if and only if ess $V_{a}^{b} f<\infty$. Moreover, the total variation of the measure $f^{\prime}$ is given by $\left\|f^{\prime}\right\|=$ ess $V_{a}^{b} f$.

Proof. First, under the assumption that ess $V_{a}^{b} f<\infty$, we will prove that $f^{\prime}$ is a measure and that

$$
\begin{equation*}
\left\|f^{\prime}\right\| \leq \operatorname{ess} V_{a}^{b} f \tag{10.18}
\end{equation*}
$$

For this purpose, choose $\varepsilon>0$ and let $f_{\varepsilon}=\psi_{\varepsilon} * f$ denote the mollifier of $f$ (see (10.9)). Consider an arbitrary partition with the property $a+\varepsilon<t_{1}<\cdots<$
$t_{m+1}<b-\varepsilon$. If a function $g$ is defined for fixed $t_{i}$ by $g(s)=f\left(t_{i}-s\right)$, then almost every $s$ is a point of approximate continuity of $g$ and therefore of $f\left(t_{i}-\cdot\right)$. Hence,

$$
\begin{aligned}
\sum_{i=1}^{m}\left|f_{\varepsilon}\left(t_{i+1}\right)-f_{\varepsilon}\left(t_{i}\right)\right| & =\sum_{i=1}^{m}\left|\int_{-\varepsilon}^{\varepsilon} \psi_{\varepsilon}(s)\left(f\left(t_{i+1}-s\right)-f\left(t_{i}-s\right)\right) d \lambda(s)\right| \\
& \leq \int_{-\varepsilon}^{\varepsilon} \psi_{\varepsilon}(s) \sum_{i=1}^{m}\left|f\left(t_{i+1}-s\right)-f\left(t_{i}-s\right)\right| d \lambda(s) \\
& \leq \int_{-\varepsilon}^{\varepsilon} \psi_{\varepsilon}(s) \operatorname{ess} V_{a}^{b} f d \lambda(s) \\
& \leq \operatorname{ess} V_{a}^{b} f
\end{aligned}
$$

In obtaining the last inequality, we have used the fact that

$$
\int_{-\varepsilon}^{\varepsilon} \psi_{\varepsilon}(s) d \lambda(s)=1
$$

Now take the supremum of the left side over all partitions and obtain

$$
\int_{a+\varepsilon}^{b-\varepsilon}\left|\left(f_{\varepsilon}\right)^{\prime}\right| d \lambda \leq \operatorname{ess} V_{a}^{b} f
$$

Let $\varphi \in C_{c}^{\infty}(a, b)$ with $|\varphi| \leq 1$. Choosing $\varepsilon>0$ such that $\operatorname{spt} \varphi \subset(a+\varepsilon, b-\varepsilon)$, we obtain

$$
\int_{a}^{b} f_{\varepsilon} \varphi^{\prime} d \lambda=-\int_{a}^{b}\left(f_{\varepsilon}\right)^{\prime} \varphi d \lambda \leq \int_{a+\varepsilon}^{b-\varepsilon}\left|\left(f_{\varepsilon}\right)^{\prime}\right| d \lambda \leq \operatorname{ess} V_{a}^{b} f
$$

Since $f_{\varepsilon}$ converges to $f$ in $L^{1}(a, b)$, it follows that

$$
\begin{equation*}
\int_{a}^{b} f \varphi^{\prime} d \lambda \leq \operatorname{ess} V_{a}^{b} f \tag{10.19}
\end{equation*}
$$

Thus, the distribution $S$ defined for all test functions $\psi$ by

$$
S(\psi)=\int_{a}^{b} f \psi^{\prime} d \lambda
$$

is a distribution of order 0 and therefore a measure since by (10.19),

$$
S(\psi)=\int_{a}^{b} f \psi^{\prime} d \lambda=\int_{a}^{b} f \varphi^{\prime}\|\psi\|_{(a, b) ; 0} d \lambda \leq \operatorname{ess} V_{a}^{b} f\|\psi\|_{(a, b) ; 0}
$$

where we have taken

$$
\varphi=\frac{\psi}{\|\psi\|_{(a, b) ; 0}}
$$

Since $S=-f^{\prime}$, we have that $f^{\prime}$ is a measure with

$$
\left\|f^{\prime}\right\|=\sup \left\{\int_{a}^{b} f \varphi^{\prime} d \lambda: \varphi \in C_{c}^{\infty}(a, b),|\varphi| \leq 1\right\} \leq \operatorname{ess} V_{a}^{b} f
$$

thus establishing (10.18) as desired.
Now for the opposite inequality. We assume that $f^{\prime}$ is a measure with finite total variation, and we will first show that $f \in L^{\infty}(a, b)$. For $0<h<(b-a) / 3$, let
$I_{h}$ denote the interval $(a+h, b-h)$ and let $\eta \in C_{c}^{\infty}\left(I_{h}\right)$ with $|\eta| \leq 1$. Then, for all sufficiently small $\varepsilon>0$, we have the mollifier $\eta_{\varepsilon}=\eta * \varphi_{\varepsilon} \in C_{c}^{\infty}(a, b)$. Thus, with the help of (10.10) and Fubini's Theorem,

$$
\begin{aligned}
\int_{I_{h}} f_{\varepsilon} \eta^{\prime} d \lambda & =\int_{a}^{b} f_{\varepsilon} \eta^{\prime} d \lambda=\int_{a}^{b}\left(f * \varphi_{\varepsilon}\right) \eta^{\prime} d \lambda \\
& =\int_{a}^{b} f\left(\varphi_{\varepsilon} * \eta\right)^{\prime} d \lambda=\int_{a}^{b} f \eta_{\varepsilon}^{\prime} d \lambda \leq\left\|f^{\prime}\right\|
\end{aligned}
$$

since $\left|\eta_{\varepsilon}\right| \leq 1$. Taking the supremum over all such $\eta$ shows that

$$
\begin{equation*}
\left\|f_{\varepsilon}^{\prime}\right\|_{L^{1} ; I_{h}} \leq\left\|f^{\prime}\right\| \tag{10.20}
\end{equation*}
$$

because

$$
\int_{I_{h}} f_{\varepsilon}^{\prime} \eta d \lambda=-\int_{I_{h}} f_{\varepsilon} \eta^{\prime} d \lambda .
$$

For arbitrary $y, z \in I_{h}$, we have

$$
f_{\varepsilon}(z)=f_{\varepsilon}(y)+\int_{y}^{z} f_{\varepsilon}^{\prime} d \lambda
$$

and therefore, taking integral averages with respect to $y$,

$$
f_{I_{h}}\left|f_{\varepsilon}(z)\right| d \lambda \leq \int_{I_{h}}\left|f_{\varepsilon}\right| d \lambda+\int_{a}^{b}\left|f_{\varepsilon}^{\prime}\right| d \lambda
$$

Thus, from Theorem 10.1 (iii) and (10.20)

$$
\left|f_{\varepsilon}\right|(z) \leq \frac{3}{b-a}\|f\|_{L^{1}(a, b)}+\left\|f^{\prime}\right\|
$$

for $z \in I_{h}$. Since $f_{\varepsilon} \rightarrow f$ a.e. in $(a, b)$ as $\varepsilon \rightarrow 0$, we conclude that $f \in L^{\infty}(a, b)$.
Now that we know that $f$ is bounded on $(a, b)$, we see that each point of approximate continuity of $f$ is also a Lebesgue point (see Exercise 4, Section 7.9). Consequently, for each partition $a<t_{1}<\cdots<t_{m+1}<b$ where each $t_{i}$ is a point of approximate continuity of $f$, reference to Theorem 10.1 (ii) yields

$$
\begin{aligned}
\sum_{i=1}^{m}\left|f\left(t_{i+1}\right)-f\left(t_{i}\right)\right| & =\lim _{\varepsilon \rightarrow 0} \sum_{i=1}^{m}\left|f_{\varepsilon}\left(t_{i+1}\right)-f_{\varepsilon}\left(t_{i}\right)\right| \\
& \leq \limsup _{\varepsilon \rightarrow 0} \int_{a}^{b}\left|f_{\varepsilon}^{\prime}\right| d \lambda \\
& \leq\|f\|(a, b)
\end{aligned}
$$

Now take the supremum over all such partitions to conclude that

$$
\operatorname{ess} V_{a}^{b} f \leq\left\|f^{\prime}\right\|
$$

## CHAPTER 11

## Functions of Several Variables

### 11.1. Differentiability


#### Abstract

Because of the central role played by absolutely continuous functions and functions of bounded variation in the development of the Fundamental Theorem of Calculus in $\mathbb{R}$, it is natural to ask whether they have analogues among functions of more than one variable. One of the main objectives of this chapter is to show that this is true. We have found that the $B V$ functions in $\mathbb{R}$ comprise a large class of functions that are differentiable almost everywhere. Although there are functions on $\mathbb{R}^{n}$ that are analogous to $B V$ functions, they are not differentiable almost everywhere. Among the functions that are often encountered in applications, Lipschitz functions on $\mathbb{R}^{n}$ form the largest class that are differentiable almost everywhere. This result is due to Rademacher and is one of the fundamental results in the analysis of functions of several variables.


The derivative of a function $f$ at a point $x_{0}$ satisfies

$$
\lim _{h \rightarrow 0}\left|\frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)-f^{\prime}\left(x_{0}\right) h}{h}\right|=0 .
$$

This limit implies that

$$
\begin{equation*}
f\left(x_{0}+h\right)-f\left(x_{0}\right)=f^{\prime}\left(x_{0}\right) h+r(h) \tag{11.1}
\end{equation*}
$$

where the "remainder" $r(h)$ is small, in the sense that

$$
\lim _{h \rightarrow 0} \frac{r(h)}{h}=0 .
$$

Note that (11.1) expresses the difference $f\left(x_{0}+h\right)-f\left(x_{0}\right)$ as the sum of the linear function that takes $h$ to $f^{\prime}\left(x_{0}\right) h$, plus a small remainder. We can therefore regard the derivative of $f$ at $x_{0}$, not as a real number, but as linear function of $\mathbb{R}$ that takes $h$ to $f^{\prime}\left(x_{0}\right) h$. Observe that every real number $\alpha$ gives rise to the linear function $L: \mathbb{R} \rightarrow \mathbb{R}, L(h)=\alpha \cdot h$. Conversely, every linear function that carries $\mathbb{R}$ to $\mathbb{R}$ is multiplication by some real number. It is this natural 1-1 correspondence which motivates the definition of differentiability for functions of several variables. Thus, if $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$, we say that $f$ is differentiable at $x_{0} \in \mathbb{R}^{n}$ provided there is a linear function $L: \mathbb{R}^{n} \rightarrow \mathbb{R}$ with the property that

$$
\begin{equation*}
\lim _{h \rightarrow 0} \frac{\left|f\left(x_{0}+h\right)-f\left(x_{0}\right)-L(h)\right|}{|h|}=0 . \tag{11.2}
\end{equation*}
$$

The linear function $L$ is called the derivative of $f$ at $x_{0}$ and is denoted by $d f\left(x_{0}\right)$. It is commonly accepted to use the term differential interchangeably with derivative.

Recall that the existence of partial derivatives at a point is not sufficient to ensure differentiability. For example, consider the following function of two variables:

$$
f(x, y)= \begin{cases}\frac{x y^{2}+x^{2} y}{x^{2}+y^{2}} & \text { if }(x, y) \neq(0,0) \\ 0 & \text { if }(x, y)=(0,0)\end{cases}
$$

Both partial derivatives are 0 at $(0,0)$, but yet the function is not differentiable there. We leave it to the reader to verify this.

We now consider Lipschitz functions $f$ on $\mathbb{R}^{n}$, see (3.10); thus there is a constant $C_{f}$ such that

$$
|f(x)-f(y)| \leq C_{f}|x-y|
$$

for all $x, y \in \mathbb{R}^{n}$. The next result is fundamental in the study of functions of several variables.
11.1. Theorem. If $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is Lipschitz, then $f$ is differentiable almost everywhere.

Proof. Step 1: Let $v \in \mathbb{R}^{n}$ with $|v|=1$. We first show that $d f_{v}(x)$ exist for $\lambda$-a.e. $x$, where $d f_{v}(x)$ denotes the directional derivative of $f$ at $x$ in the direction of $v$.

Let

$$
N_{v}:=\mathbb{R}^{n} \cap\left\{x: d f_{v}(x) \text { fails to exist }\right\}
$$

For each $x \in \mathbb{R}^{n}$ define

$$
\overline{d f_{v}}(x):=\limsup _{t \rightarrow 0} \frac{f(x+t v)-f(x)}{t}
$$

and

$$
\underline{d f_{v}}(x):=\liminf _{t \rightarrow 0} \frac{f(x+t v)-f(x)}{t}
$$

Notice that

$$
\begin{equation*}
N_{v}=\left\{x \in \mathbb{R}^{n}: \underline{d f_{v}}(x)<\overline{d f_{v}}(x)\right\} . \tag{11.3}
\end{equation*}
$$

From Definition 3.63, it is clear that

$$
\overline{d f_{v}}(x)=\lim _{k \rightarrow \infty}\left(\sup _{\substack{0<|t|<1 / k \\ t \in \mathbb{Q}, k \in \mathbb{N}}} \frac{f(x+t v)-f(x)}{t}\right)
$$

We now claim that

$$
x \rightarrow \overline{d f_{v}}(x) \text { is a Borel measurable function of } x
$$

Indeed, the rational numbers $0<|t|<\frac{1}{k}$ can be enumerated as $t_{1}^{k}, t_{2}^{k}, t_{3}^{k}, \ldots$. If we define for each $i=1,2, \ldots$ and fixed $k$ the sequence $G_{i}^{k}(x):=\frac{f\left(x+t_{i}^{k} v\right)-f(x)}{t_{i}^{k}}$ then, since $f$ is continuous, the functions $G_{i}^{k}$ are Borel measurable and hence $F^{k}(x):=$ $\sup _{i}\left\{G_{i}^{k}(x)\right\}$ is also Borel measurable. Note that

$$
F^{k}(x)=\sup _{\substack{0<|t|<1 / k \\ t \in \mathbb{Q}}} \frac{f(x+t v)-f(x)}{t}
$$

Since the pointwise limit of measurable functions is again measurable it follows that

$$
\begin{equation*}
\overline{d f_{v}}(x)=\lim _{k \rightarrow \infty} F^{k}(x) \tag{11.4}
\end{equation*}
$$

is Borel measurable. Proceeding as before we also have

$$
x \rightarrow \underline{d f_{v}}(x) \text { is a Borel measurable function of } x
$$

Therefore, from (11.3) we conclude

$$
\begin{equation*}
N_{v} \text { is a Borel set. } \tag{11.5}
\end{equation*}
$$

Now we proceed to show that

$$
\begin{equation*}
H^{1}\left(N_{v} \cap L\right)=0, \text { for each line } L \text { parallel to } v \tag{11.6}
\end{equation*}
$$

In order to prove (11.6) we consider the line $L_{x}$ that contains $x \in \mathbb{R}^{n}$ and is parallel to $v$. We now consider the restriction of $f$ to $L_{x}$ given by

$$
\gamma: \mathbb{R} \rightarrow \mathbb{R}, \gamma(t)=f(x+t v)
$$

Note that $\gamma$ is Lipschitz in $\mathbb{R}$ since $f$ is Lipschitz in $\mathbb{R}^{n}$. Therefore, $\gamma$ is absolutely continuous and hence differentiable at $\lambda_{1}$-a.e. $t$. Let $A_{v}=\{t \in \mathbb{R}: \gamma(t)$ is not differentiable $\}$ and $R_{v}=z\left(A_{v}\right)$, where $z: \mathbb{R} \rightarrow \mathbb{R}^{n}$ is given by $z(t)=x+t v$. Since $z$ is Lipschitz (Exercise 10, Section 4.7) implies that $H^{1}\left(R_{v}\right) \leq C H^{1}\left(A_{v}\right)=0$, which gives

$$
\begin{equation*}
H^{1}\left(R_{v}\right)=0 \tag{11.7}
\end{equation*}
$$

We have that $t_{0} \notin A_{v}$ if and only if $z\left(t_{0}\right)=x+t_{0} v \notin N_{v}$. Indeed, for such $t_{0} \gamma^{\prime}\left(t_{0}\right)$ exists and:

$$
\begin{aligned}
\gamma^{\prime}\left(t_{0}\right) & =\lim _{t \rightarrow 0} \frac{\gamma(t)-\gamma\left(t_{0}\right)}{t-t_{0}}=\lim _{t \rightarrow 0} \frac{f(x+t v)-f\left(x+t_{0} v\right)}{t-t_{0}} \\
& =\lim _{t \rightarrow 0} \frac{f\left(x+t_{0} v+\left(t-t_{0}\right) v\right)-f\left(x+t_{0} v\right)}{t-t_{0}} \\
& =\lim _{h \rightarrow 0} \frac{f\left(x+t_{0} v+h v\right)-f\left(x+t_{0} v\right)}{h} ; \text { with } h=t-t_{0} \\
& =\lim _{h \rightarrow 0} \frac{f\left(z\left(t_{0}\right)+h v\right)-f\left(z\left(t_{0}\right)\right)}{h}=d f_{v}\left(z\left(t_{0}\right)\right) .
\end{aligned}
$$

Hence the directional derivative of $f$ exists at the point $z\left(t_{0}\right)=x+t_{0} v$. It is now clear that

$$
R_{v}=N_{v} \cap L_{x}
$$

and from (11.7) we conclude

$$
\begin{equation*}
H^{1}\left(N_{v} \cap L_{x}\right)=0 \tag{11.8}
\end{equation*}
$$

Since (11.8) holds for every arbitrary $x$, we have

$$
\begin{equation*}
H^{1}\left(N_{v} \cap L\right)=0 \text { for every line parallel to } v \tag{11.9}
\end{equation*}
$$

Therefore, since $N_{v}$ is a Borel set, we can appeal to Fubini's Theorem to conclude that

$$
\begin{equation*}
\lambda\left(N_{v}\right)=0 . \tag{11.10}
\end{equation*}
$$

Step 2: Our next objective is to prove that

$$
\begin{equation*}
d f_{v}(x)=\nabla f(x) \cdot v \tag{11.11}
\end{equation*}
$$

for $\lambda$-almost all $x \in \mathbb{R}^{n}$. Here, $\nabla f$ denotes the gradient of $f$ as defined in (10.6). Of course, this formula is valid for all $x$ if $f \in C^{\infty}$. As a result of (11.10), we see that $\nabla f(x)$ exists for $\lambda$-almost all $x$. To establish (11.11), we begin with an observation that follows directly from Theorem 11.4 (change of variables formula), which will be established later:

$$
\int_{\mathbb{R}^{n}}\left(\frac{f(x+t v)-f(x)}{t}\right) \varphi(x) d \lambda(x)=
$$

$$
\begin{equation*}
-\int_{\mathbb{R}^{n}} f(x)\left(\frac{\varphi(x)-\varphi(x-t v)}{t}\right) d \lambda(x) \tag{11.12}
\end{equation*}
$$

whenever $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. Letting $t$ assume the values $t=1 / k$ for all nonnegative integers $k$, we have

$$
\begin{equation*}
\left|\frac{f\left(x+\frac{1}{k} v\right)-f(x)}{\frac{1}{k}}\right| \leq C_{f}|v|=C_{f} . \tag{11.13}
\end{equation*}
$$

From (11.12),(11.13) and the Dominated Convergence Theorem we have

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} d f_{v}(x) \varphi(x) d \lambda(x) & =\int_{\mathbb{R}^{n}} \lim _{k \rightarrow \infty} \frac{f\left(x+\frac{1}{k} v\right)-f(x)}{\frac{1}{k}} \varphi(x) d \lambda(x) \\
& =\lim _{k \rightarrow \infty} \int_{\mathbb{R}^{n}} \frac{f\left(x+\frac{1}{k} v\right)-f(x)}{\frac{1}{k}} \varphi(x) d \lambda(x) \\
& =\lim _{k \rightarrow \infty} \int_{\mathbb{R}^{n}} \frac{\varphi\left(x-\frac{1}{k} v\right)-\varphi(x)}{\frac{1}{k}} f(x) d \lambda(x) \\
& =-\int_{\mathbb{R}^{n}} \lim _{k \rightarrow \infty} \frac{\varphi\left(x-\frac{1}{k} v\right)-\varphi(x)}{\frac{-1}{k}} f(x) d \lambda(x) \\
& =-\int_{\mathbb{R}^{n}} d \varphi_{v}(x) f(x) d \lambda(x)
\end{aligned}
$$

We recall that $f$ is absolutely continuous on lines and therefore, using Fubini's Theorem and integration by parts along lines we compute

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} d f_{v}(x) \varphi(x) d \lambda(x) & =-\int_{\mathbb{R}^{n}} d \varphi_{v}(x) f(x) d \lambda(x) \\
& =-\int_{\mathbb{R}^{n}} f(x) \nabla \varphi(x) \cdot v d \lambda(x) \\
& =-\sum_{i=1}^{n} v_{i} \int_{\mathbb{R}^{n}} f(x) \frac{\partial \varphi}{\partial x_{i}}(x) d \lambda(x) \\
& =\sum_{i=1}^{n} v_{i} \int_{\mathbb{R}^{n}} \frac{\partial f}{\partial x_{i}}(x) \varphi(x) d \lambda(x) \\
& =\int_{\mathbb{R}^{n}} \nabla f(x) \cdot v \quad \varphi(x) d \lambda(x) .
\end{aligned}
$$

Thus

$$
\int_{\mathbb{R}^{n}} d f_{v}(x) \varphi(x) d \lambda(x)=\int_{\mathbb{R}^{n}} \nabla f(x) \cdot v \varphi(x) d \lambda(x), \text { for all } \varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)
$$

Hence

$$
d f_{v}(x)=\nabla f(x) \cdot v, \lambda \text {-a.e. } x .
$$

Now chose $\left\{v_{k}\right\}_{k=1}^{\infty}$ be a countable dense subset of $\partial B(0,1)$. Observe that there is a set $E, \lambda(E)=0$, such that

$$
d f_{v}(x)=\nabla f(x) \cdot v_{k}, \text { for all } x \in \mathbb{R}^{n} \backslash E .
$$

Step 3: We will now show that $f$ is differentiable at each point $x \in \mathbb{R}^{n} \backslash E$. For $x \in \mathbb{R}^{n} \backslash E$ and $v \in \partial B(0,1)$ we define

$$
Q(x, v, t):=\frac{f(x+t v)-f(x)}{t}-\nabla f(x) \cdot v, \quad t \neq 0
$$

For $v, v^{\prime} \in \partial B(0,1)$ note that:

$$
\begin{aligned}
\left|Q(x, v, t)-Q\left(x, v^{\prime}, t\right)\right| & \leq \frac{\left|f(x+t v)-f\left(x+t v^{\prime}\right)\right|}{|t|}+\left|\nabla f(x) \cdot\left(v-v^{\prime}\right)\right| \\
& \leq C_{f}\left|v-v^{\prime}\right|+n C_{f}\left|v-v^{\prime}\right| \\
& =C_{f}(n+1)\left|v-v^{\prime}\right|
\end{aligned}
$$

Hence

$$
\begin{equation*}
\left|Q(x, v, t)-Q\left(x, v^{\prime}, t\right)\right| \leq C_{f}(n+1)\left|v-v^{\prime}\right|, \quad v, v^{\prime} \in \partial B(0,1) \tag{11.14}
\end{equation*}
$$

Since $\left\{v_{k}\right\}$ is dense in the compact set $\partial B(0,1)$, it follows that for every $\varepsilon>0$ there exists $N$ sufficiently large such that, for every $v \in \partial B(0,1)$,

$$
\begin{equation*}
\left|v-v_{k}\right|<\frac{\varepsilon}{2(n+1) C_{f}}, \text { for some } k \in\{1, \ldots, N\} \tag{11.15}
\end{equation*}
$$

Thus, for every $v \in \partial B(0,1)$, there exists $k \in\{1, \ldots, N\}$ such that

$$
\begin{equation*}
|Q(x, v, t)| \leq\left|Q\left(x, v_{k}, t\right)\right|+\left|Q(x, v, t)-Q\left(x, v_{k}, t\right)\right| \tag{11.16}
\end{equation*}
$$

Since $d f_{v_{k}}(x)=\nabla f(x) \cdot v_{k}$, then for $0<|t| \leq \delta$ we have $\left\lvert\, Q\left(x, v_{k}, t \left\lvert\, \leq \frac{\varepsilon}{2}\right.\right.$. Thus from \right. (11.14), (11.15), (11.16)

$$
|Q(x, v, t)| \leq \frac{\varepsilon}{2}+\frac{\varepsilon}{2}=\varepsilon, \quad 0 \leq|t| \leq \delta
$$

We have shown that

$$
\lim _{t \rightarrow 0}|Q(x, v, t)|=0
$$

which is

$$
\begin{equation*}
\lim _{t \rightarrow 0}\left|\frac{f(x+t v)-f(x)}{t}-\nabla f(x) \cdot v\right|=0 . \tag{11.17}
\end{equation*}
$$

The last step is to show that (11.17) implies that $f$ is differentiable at every $x \in$ $\mathbb{R}^{n} \backslash E$. Choose any $y \in \mathbb{R}^{n}, y \neq x$. Let

$$
v:=\frac{y-x}{\|y-x\|}, \quad y=x+t v, \quad t=|y-x|
$$

From (11.17)

$$
\begin{equation*}
\lim _{|y-x| \rightarrow 0} \frac{|f(y)-f(x)-\nabla f(x) \cdot(y-x)|}{|y-x|}=0 \tag{11.18}
\end{equation*}
$$

or, with $h:=y-x$,

$$
\begin{equation*}
\lim _{|h| \rightarrow 0} \frac{|f(x+h)-f(x)-\nabla f(x) \cdot h|}{|h|}=0 \tag{11.19}
\end{equation*}
$$

which means the $f$ is differentiable at $x$. Since $x \notin E$ we conclude $f$ is differentiable almost everywhere.

The concept of differentiability for a transformation $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is virtually the same as in (11.2). Thus, we say that $T$ is differentiable at $x_{0} \in \mathbb{R}^{n}$ if there is a linear mapping $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ such that

$$
\begin{equation*}
\lim _{h \rightarrow 0} \frac{\left|T\left(x_{0}+h\right)-T\left(x_{0}\right)-L(h)\right|}{|h|}=0 . \tag{11.20}
\end{equation*}
$$

As in the case when $T$ is real-valued, we call $L$ the derivative of $T$ at $x_{0}$. We will denote the linear function $L$ by $d T\left(x_{0}\right)$. Thus, $d T\left(x_{0}\right)$ is a linear transformation, and when it is applied to a vector $v \in \mathbb{R}^{n}$, we will write $d T\left(x_{0}\right)(v)$. Writing $T$ in terms of its coordinate functions, $T=\left(T^{1}, T^{2}, \ldots, T^{m}\right)$, it is easy to see that $T$ is differentiable at a point $x_{0}$ if and only if each $T^{i}$ is. Consequently, the following corollary is immediate.
11.2. Corollary. If $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is a Lipschitz transformation, then $T$ is differentiable $\lambda$-almost everywhere.

## Exercises for Section 11.1

1. Prove that a linear map $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is Lipschitz.
2. Show that a linear map $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ satisfies condition $N$ and therefore leaves Lebesgue measurable sets invariant.
3. Use Fubini's Theorem to prove

$$
\frac{\partial^{2} f}{\partial x \partial y}=\frac{\partial^{2} f}{\partial y \partial x}
$$

if $f \in C^{2}\left(\mathbb{R}^{2}\right)$. Use this result to conclude that all second order mixed partials are equal if $f \in C^{2}\left(\mathbb{R}^{n}\right)$.
4. Let $C^{1}[0,1]$ denote the space of functions on $[0,1]$ that have continuous derivatives on $[0,1]$, including one-sided derivatives at the endpoints. Define a norm on $C^{1}[0,1]$ by

$$
\|f\|:=\sup _{x \in[0,1]}|f(x)|+\sup _{x \in[0,1]} f^{\prime}(x) .
$$

Prove that $C^{1}[0,1]$ with this norm is a Banach space.

### 11.2. Change of Variable

We give a treatment of the behavior of the integral when the integrand is subjected to a change of variables by a Lipschitz transformation $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$.

Consider $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ with $T=\left(T^{1}, T^{2}, \ldots, T^{n}\right)$. If $T$ is differentiable at $x_{0}$, then it follows immediately from definitions that each of the partial derivatives

$$
\frac{\partial T^{i}}{\partial x_{j}} i, j=1,2, \ldots, n
$$

exists at $x_{0}$. The linear mapping $L=d T\left(x_{0}\right)$ in (11.20) can be represented by the $n \times n$ matrix

$$
d T\left(x_{0}\right)=\left[\begin{array}{ccc}
\frac{\partial T^{1}}{\partial x_{1}} & \cdots & \frac{\partial T^{1}}{\partial x_{n}} \\
\vdots & & \vdots \\
\frac{\partial T^{n}}{\partial x_{1}} & \cdots & \frac{\partial T^{n}}{\partial x_{n}}
\end{array}\right]
$$

where it is understood that each partial derivative is evaluated at $x_{0}$. The determinant of $\left[d T\left(x_{0}\right)\right]$ is called the Jacobian of $T$ at $x_{0}$ and is denoted by $J T\left(x_{0}\right)$.

Recall from elementary linear algebra that a linear map $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ can be identified with the $n \times n$ matrix $\left(L_{i j}\right)$ where

$$
L_{i j}=L\left(e_{i}\right) \cdot e_{j}, \quad i, j=1,2, \ldots, n
$$

and where $\left\{e_{i}\right\}$ is the standard basis for $\mathbb{R}^{n}$. The determinant of $L_{i j}$ is denoted by $\operatorname{det} L$. If $M: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is also a linear map,

$$
\begin{equation*}
\operatorname{det}(M \circ L)=(\operatorname{det} M) \cdot(\operatorname{det} L) \tag{11.21}
\end{equation*}
$$

Every nonsingular $n \times n$ matrix $\left(L_{i j}\right)$ can be row-reduced to the identity matrix. That is, $L$ can be written as the composition of finitely many linear transformations of the following three types:
(i) $L_{1}\left(x_{1}, \ldots, x_{i}, \ldots, x_{n}\right)=\left(x_{1}, \ldots, c x_{i}, \ldots, x_{n}\right), \quad 1 \leq i \leq n, c \neq 0$.
(ii) $L_{2}\left(x_{1}, \ldots, x_{i}, \ldots, x_{n}\right)=\left(x_{1}, \ldots, x_{i}+c x_{k}, \ldots, x_{n}\right), \quad 1 \leq i \leq n, k \neq i, c \neq 0$.
(iii) $L_{3}\left(x_{1}, \ldots, x_{i}, \ldots, x_{j}, \ldots, x_{n}\right)=\left(x_{1}, \ldots, x_{j}, \ldots, x_{i}, \ldots, x_{n}\right), \quad 1 \leq i<j \leq n$.

This leads to the following geometric interpretation of the determinant.
11.3. THEOREM. If $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a nonsingular linear map, then $L(E)$ is Lebesgue measurable whenever $E$ is Lebesgue measurable and

$$
\begin{equation*}
\lambda[L(E)]=|\operatorname{det} L| \lambda(E) \tag{11.22}
\end{equation*}
$$

Proof. Exercise 2, Section 11.1, implies that $L(E)$ is Lebesgue measurable. In view of (11.21), it suffices to prove (11.22) when $L$ is of the three types mentioned above. In the case of $L_{3}$, we use Fubini's Theorem and interchange the order of integration. In the case of $L_{1}$ or $L_{2}$, we integrate first respect to $x_{i}$ to arrive at the formulas of the form

$$
\begin{aligned}
|c| \lambda_{1}(A) & =\lambda_{1}(c A) \\
\lambda_{1}(A) & =\lambda_{1}(c+A)
\end{aligned}
$$

11.4. Theorem. If $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a nonsingular linear mapping, then

$$
\int_{\mathbb{R}^{n}} f \circ L|\operatorname{det} L| d \lambda=\int_{\mathbb{R}^{n}} f d \lambda .
$$

Proof. Assume first that $f \geq 0$. Note that

$$
\{f>t\}=L(\{f \circ L>t\})
$$

and therefore

$$
\lambda(\{f>t\})=\lambda(L(\{f \circ L>t\}))=|\operatorname{det} L| \lambda(\{f \circ L>t\})
$$

for $t \geq 0$. Now apply Theorem 6.60 to obtain our desired result. The general case follows by writing $f=f^{+}-f^{-}$.

Our next result deals with an approximation property of Lipschitz transformations with nonvanishing Jacobians. First, we recall that a linear transformation $L: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ can be identified with its $n \times n$ matrix. Let $\mathcal{R}$ denote the family of $n \times n$ matrices whose entries are rational numbers. Clearly, $\mathcal{R}$ is countable. Furthermore, given an arbitrary nonsingular linear transformation $L$ and $\varepsilon>0$, there exists $R \in \mathcal{R}$ such that

$$
\begin{array}{r}
|L(x)-R(x)|<\varepsilon \\
\left|L^{-1}(x)-R^{-1}(x)\right|<\varepsilon \tag{11.23}
\end{array}
$$

whenever $|x| \leq 1$. Using linearity, this implies

$$
\begin{aligned}
|L(x)-R(x)| & <\varepsilon|x| \\
\left|L^{-1}(x)-R^{-1}(x)\right| & <\varepsilon|x|
\end{aligned}
$$

for all $x \in \mathbb{R}^{n}$. Also, (11.23) implies

$$
\left|L \circ R^{-1}(x-y)\right| \leq(1+\varepsilon)|x-y|
$$

and

$$
\left|R \circ L^{-1}(x-y)\right| \leq(1+\varepsilon)|x-y|
$$

for all $x, y \in \mathbb{R}^{n}$. That is, the Lipschitz constants (see (3.10)) of $L \circ R^{-1}$ and $R \circ L^{-1}$ satisfy

$$
\begin{equation*}
C_{L \circ R^{-1}}<1+\varepsilon \quad \text { and } \quad C_{R \circ L^{-1}}<1+\varepsilon . \tag{11.24}
\end{equation*}
$$

11.5. Theorem. Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a continuous transformation and set

$$
B=\{x: d T(x) \text { exists, } J T(x) \neq 0\} .
$$

Given $t>1$ there exists a countable collection of Borel sets $\left\{B_{k}\right\}_{k=1}^{\infty}$ such that
(i) $B=\bigcup_{k=1}^{\infty} B_{k}$,
(ii) The restriction of $T$ to $B_{k}$ (denoted by $T_{k}$ ) is univalent,
(iii) For each positive integer $k$, there exists a nonsingular linear transformation $L_{k}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ such that the Lipschitz constants of $T_{k} \circ L_{k}^{-1}$ and $L_{k} \circ T_{k}^{-1}$ satisfy

$$
C_{T_{k} \circ L_{k}^{-1}} \leq t, C_{L_{k} \circ T_{k}^{-1}} \leq t
$$

and

$$
t^{-n}\left|\operatorname{det} L_{k}\right| \leq|J T(x)| \leq t^{n}\left|\operatorname{det} L_{k}\right|
$$

for all $x \in B_{k}$.
Proof. For fixed $t>1$ choose $\varepsilon>0$ such that

$$
\frac{1}{t}+\varepsilon<1<t-\varepsilon
$$

Let $C$ be a countable dense subset of $B$ and let $\mathcal{R}$ (as introduced above) be the family of linear transformations whose matrices have rational entries. Now, for each $c \in C, R \in \mathcal{R}$ and each positive integer $i$, define $E(c, R, i)$ to be the set of all $b \in B \cap B(c, 1 / i)$ that satisfy

$$
\begin{equation*}
\left(\frac{1}{t}+\varepsilon\right)|R(v)| \leq|d T(b)(v)| \leq(t-\varepsilon)|R(v)| \tag{11.25}
\end{equation*}
$$

for all $v \in \mathbb{R}^{n}$ and

$$
\begin{equation*}
|T(a)-T(b)-d T(b)(a-b)| \leq \varepsilon|R(a-b)| \tag{11.26}
\end{equation*}
$$

for all $a \in B(b, 2 / i)$. Since $T$ is continuous, each partial derivative of each coordinate function of $T$ is a Borel function. Thus, it is an easy exercise to prove that each $E(c, R, i)$ is a Borel set. Observe that (11.25) and (11.26) imply

$$
\begin{equation*}
\frac{1}{t}|R(a-b)| \leq|T(a)-T(b)| \leq t|R(a-b)| \tag{11.27}
\end{equation*}
$$

for all $b \in E(c, R, i)$ and $a \in B(b, 2 / i)$.
Next, we will show that

$$
\begin{equation*}
\left(\frac{1}{t}+\varepsilon\right)^{n}|\operatorname{det} R| \leq|J T(b)| \leq(t-\varepsilon)^{n}|\operatorname{det} R| \tag{11.28}
\end{equation*}
$$

for $b \in E(c, R, i)$. For the proof of this, let $L=d T(b)$. Then, from (11.25) we have

$$
\begin{equation*}
\left(\frac{1}{t}+\varepsilon\right)|R(v)| \leq|L(v)| \leq(t-\varepsilon)|R(v)| \tag{11.29}
\end{equation*}
$$

whenever $v \in \mathbb{R}^{n}$ and therefore,

$$
\begin{equation*}
\left(\frac{1}{t}+\varepsilon\right)|v| \leq\left|L \circ R^{-1}(v)\right| \leq(t-\varepsilon)|v| \tag{11.30}
\end{equation*}
$$

for $v \in \mathbb{R}^{n}$. This implies that

$$
L \circ R^{-1}[B(0,1)] \subset B(0, t-\varepsilon)
$$

and reference to Theorem 11.3 yields

$$
\left|\operatorname{det}\left(L \circ R^{-1}\right)\right| \alpha(n) \leq \lambda[B(0, t-\varepsilon)]=\alpha(n)(t-\varepsilon)^{n}
$$

where $\alpha(n)$ denotes the volume of the unit ball. Thus,

$$
|\operatorname{det} L| \leq(t-\varepsilon)^{n}|\operatorname{det} R| .
$$

This proves one part of (11.28). The proof of the other part is similar.
We now are ready to define the Borel sets $B_{k}$ appearing in the statement of our Theorem. Since the parameters $c, R$, and $i$ used to define the sets $E(c, R, i)$ range over countable sets, the collection $\{E(c, R, i)\}$ is itself countable. We will relable the sets $E(c, R, i)$ as $\left\{B_{k}\right\}_{k=1}^{\infty}$.

To show that property (i) holds, choose $b \in B$ and let $L=d T(b)$ as above. Now refer to (11.24) to find $R \in \mathcal{R}$ such that

$$
C_{L \circ R^{-1}}<\left(\frac{1}{t}+\varepsilon\right)^{-1} \quad \text { and } \quad C_{R \circ L^{-1}}<t-\varepsilon
$$

Using the definition of the differentiability of $T$ at $b$, select a nonnegative integer $i$ such that

$$
|T(a)-T(b)-d T(b) \cdot(a-b)| \leq \frac{\varepsilon}{C_{R^{-1}}}|a-b| \leq \varepsilon|R(a-b)|
$$

for all $a \in B(b, 2 / i)$. Now choose $c \in C$ such that $|b-c|<1 / i$ and conclude that $b \in E(c, R, i)$. Since this holds for all $b \in B$, property (i) holds.

To prove (ii), choose any set $B_{k}$. It is one of the sets of the form $E(c, R, i)$ for some $c \in C, R \in \mathcal{R}$, and some nonnegative integer $i$. According to (11.27),

$$
\frac{1}{t}|R(a-b)| \leq|T(a)-T(b)| \leq t|R(a-b)|
$$

for all $b \in B_{k}, a \in B(b, 2 / i)$. Since $B_{k} \subset B(c, 1 / i) \subset B(b, 2 / i)$, we thus have

$$
\begin{equation*}
\frac{1}{t}|R(a-b)| \leq|T(a)-T(b)| \leq t|R(a-b)| \tag{11.31}
\end{equation*}
$$

for all $a, b \in B_{k}$. Hence, $T$ restricted to $B_{k}$ is univalent.
With $B_{k}$ of the form $E(c, R, i)$ as in the preceding paragraph, we define $L_{k}=R$. The proof of (iii) follows from (11.31) and (11.28). They imply

$$
C_{T_{k} \circ L_{k}^{-1}} \leq t, \quad C_{L_{k} \circ T_{k}^{-1}} \leq t,
$$

and

$$
t^{-n}\left|\operatorname{det} L_{k}\right| \leq\left|J T_{k}\right| \leq t^{n}\left|\operatorname{det} L_{k}\right|
$$

since $\varepsilon$ is arbitrary.
We now proceed to develop the analog of Banach's Theorem (Theorem 7.33) for Lipschitz mappings $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$. As in (7.40), for $E \subset \mathbb{R}^{n}$ we define

$$
N(T, E, y)
$$

as the (possibly infinite) number of points in $E \cap T^{-1}(y)$.
11.6. Lemma. Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a Lipschitz transformation. If

$$
E:=\left\{x \in \mathbb{R}^{n}: T \text { is differentiable at } x \text { and } J T(x)=0\right\}
$$

then

$$
\lambda(T(E))=0
$$

Proof. By Rademacher's theorem, we know that $T$ is differentiable almost everywhere. Since each entry of $d T$ is a measurable function, it follows that the set $E$ is measurable and therefore, so is $T(E)$ since $T$ preserves sets of measure zero. It is sufficient to prove that $T\left(E_{R}\right)$ has measure zero for each $R>0$ where $E_{R}:=E \cap B(0, R)$. Let $\varepsilon \in(0,1)$ and fix $x \in E_{R}$. Since $T$ is differentiable at $x$, , there exists $0<\delta_{x}<1$ such that

$$
\begin{equation*}
|T(y)-T(x)-L(y-x)| \leq \varepsilon|y-x| \text { for all } y \in B\left(x, \delta_{x}\right) \tag{11.32}
\end{equation*}
$$

where $L:=d T(x)$. We know that $J T(x)=0$, and so $L$ is represented by a singular matrix. Therefore, there is a linear subspace $H$ of $\mathbb{R}^{n}$ of dimension $n-1$ such that $L$ maps $\mathbb{R}^{n}$ into $H$. From (11.32) and the triangle inequality, we have

$$
|T(y)-T(x)| \leq|L(y-x)|+\varepsilon|y-x|
$$

Let $\kappa_{L}$ denote the Lipschitz constant of $L:|L(y)| \leq \kappa_{L}|y|$ for all $y \in \mathbb{R}^{n}$. Also, let $\kappa_{T}$ denote the Lipschitz constant of $T$. Hence, for $x \in E_{R}, 0<r<\delta_{x}<1$ and $y \in B(x, r)$ it follows that

$$
\begin{aligned}
|L(y)-L(x)| & \leq \kappa_{L}|y-x| \leq \kappa_{L}[|y|+R] \\
& =\kappa_{L}[|y-x+x|+R] \leq \kappa_{L}[r+2 R]
\end{aligned}
$$

This implies

$$
|L(y)| \leq|L(x)|+\kappa_{L}[r+2 R] \leq \kappa_{L} R+\kappa_{L}[r+2 R]=\kappa_{L}[r+3 R]
$$

With $v:=T(x)-L(x)$ we see that $|v| \leq|T(x)|+|L(x)| \leq \kappa_{L}|x|+\kappa_{L}|x| \leq$ $R\left(\kappa_{T}+\kappa_{L}\right)$ and that (11.32) becomes

$$
|T(y)-L(y)-v| \leq \varepsilon|y-x| ;
$$

that is, $T(y)$ is within distance $\varepsilon r$ of $L(y)$ translated by the vector $v$. This means that the set $T(B(x, r))$ is contained within an $\varepsilon r$-neighborhood of a bounded set in an affine space of dimension $n-1$. The bound on this set depends only on $r, R, \kappa_{L}$ and $\kappa_{T}$. Thus there is a constant $M=M\left(R, \kappa_{L}, \kappa_{T}\right)$ such that

$$
\lambda(T(B(x, r))) \leq \varepsilon r M r^{n-1}
$$

The collection of balls $\{B(x, r)\}$ with $x \in E_{R}$ and $0<r<\delta_{x}$ determines a Vitali covering of $E_{R}$ and accordingly there is a countable disjoint subcollection $B_{i}:=$ $B\left(x_{i}, r_{i}\right)$ whose union contains almost all of $E_{R}$ :

$$
\lambda(F)=0 \text { where } F:=E_{R} \backslash \bigcup_{i=1}^{\infty} B_{i} .
$$

Since $T$ is Lipschitz we know that $\lambda(T(F))=0$. Therefore, because

$$
E_{R} \subset F \cup \bigcup_{i=1}^{\infty} B_{i},
$$

we have

$$
T\left(E_{R}\right) \subset T(F) \cup \bigcup_{i=1}^{\infty} T\left(B_{i}\right)
$$

and thus,

$$
\begin{aligned}
\lambda\left(T\left(E_{R}\right)\right) & \leq \sum_{i=1}^{\infty} T\left(B_{i}\right) \\
& \leq \sum_{i=1}^{\infty} \varepsilon r_{i} M r_{i}^{n-1} \\
& \leq \varepsilon M \sum_{i=1}^{\infty} \lambda\left(B\left(x_{i}, r_{i}\right)\right) .
\end{aligned}
$$

Since the balls $B\left(x_{i}, r_{i}\right)$ are disjoint and all contained in $B(0, R+1)$ and since $\varepsilon>0$ was chosen arbitrarily, we conclude that $\lambda\left(T\left(E_{R}\right)\right)=0$, as desired.
11.7. Theorem. (Area Formula) Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be Lipschitz. Then for each Lebesgue measurable set $E \subset \mathbb{R}^{n}$,

$$
\int_{E}|J T(x)| d \lambda(x)=\int_{\mathbb{R}^{n}} N(T, E, y) d \lambda(y)
$$

Proof. Since $T$ is Lipschitz, we know that $T$ carries sets of measure zero into sets of measure zero. Thus, by Rademacher's Theorem, we might as well assume that $d T(x)$ exists for all $x \in E$. Furthermore, it is easy to see that we may assume $\lambda(E)<\infty$.

In view of Lemma 11.6 it suffices to treat the case when $|J T| \neq 0$ on $E$. Fix $t>1$ and let $\left\{B_{k}\right\}$ denote the sets provided by Theorem 11.5. By Lemma 4.7, we may assume that the sets $\left\{B_{k}\right\}$ are disjoint. For the moment, select some set $B_{k}$ and set $B=B_{k}$. For each positive integer $j$, let $\mathcal{Q}_{j}$ denote a decomposition of $\mathbb{R}^{n}$ into disjoint, "half-open" cubes with side length $1 / j$ and of the form $\left[a_{1}, b_{1}\right) \times \cdots \times$ $\left[a_{n}, b_{n}\right)$. Set

$$
F_{i, j}=B \cap E \cap Q_{i}, \quad Q_{i} \in \mathcal{Q}_{j}
$$

Then, for fixed $j$ the sets $F_{i, j}$ are disjoint and

$$
B \cap E=\bigcup_{i=1}^{\infty} F_{i, j} .
$$

Furthermore, the sets $T\left(F_{i, j}\right)$ are measurable since $T$ is Lipschitz. Thus, the functions $g_{j}$ defined by

$$
g_{j}=\sum_{i=1}^{\infty} \chi_{T\left(F_{i, j}\right)}
$$

are measurable. Now $g_{j}(y)$ is the number of sets $\left\{F_{i, j}\right\}$ such that $F_{i, j} \cap T^{-1}(y) \neq 0$ and

$$
\lim _{j \rightarrow \infty} g_{j}(y)=N(T, B \cap E, y)
$$

An application of the Monotone Convergence Theorem yields

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \sum_{i=1}^{\infty} \lambda\left[T\left(F_{i, j}\right)\right]=\int_{\mathbb{R}^{n}} N(T, B \cap E, y) d \lambda(y) \tag{11.33}
\end{equation*}
$$

Let $L_{k}$ and $T_{k}$ be as in Theorem 11.5. Then, recalling that $B=B_{k}$, we obtain

$$
\begin{align*}
& \begin{array}{l}
\lambda\left[T\left(F_{i, j}\right)\right]=\lambda\left[T_{k}\left(F_{i, j}\right)\right] \\
\quad=\lambda\left[\left(T_{k} \circ L_{k}^{-1} \circ L_{k}\right)\left(F_{i, j}\right)\right] \leq t^{n} \lambda\left[L_{k}\left(F_{i, j}\right)\right] \\
\quad \lambda\left[L_{k}\left(F_{i, j}\right)\right]
\end{array} \quad \begin{array}{l}
\quad \lambda\left[\left(L_{k} \circ T_{k}^{-1} \circ T_{k}\right)\left(F_{i, j}\right)\right] \\
\quad \leq t^{n} \lambda\left[T_{k}\left(F_{i, j}\right)\right]=t^{n} \lambda\left[T\left(F_{i, j}\right)\right]
\end{array} \tag{11.34}
\end{align*}
$$

and thus

$$
\begin{array}{rlr}
t^{-2 n} \lambda\left[T\left(F_{i, j}\right)\right] & \leq t^{-n} \lambda\left[L_{k}\left(F_{i, j}\right)\right] & \text { by (11.34) } \\
& =t^{-n}\left|\operatorname{det} L_{k}\right| \lambda\left(F_{i, j}\right) & \\
& \text { by Theorem } 11.3 \\
& \leq \int_{F_{i, j}}|J T| d \lambda & \\
& \leq t^{n}\left|\operatorname{let} L_{k}\right| \lambda\left(F_{i, j}\right) & \\
& \text { by Theorem } 11.5 \\
& =t^{n} \lambda\left[L_{k}\left(F_{i, j}\right)\right] & \\
& \leq t^{2 n} \lambda\left[T\left(F_{i, j}\right)\right] . & \text { by Theorem } 11.5 \\
& & \text { by }(11.35)
\end{array}
$$

With $j$ fixed, sum on $i$ and use the fact that the sets $F_{i, j}$ are disjoint:

$$
\begin{equation*}
t^{-2 n} \sum_{i=1}^{\infty} \lambda\left[T\left(F_{i, j}\right)\right] \leq \int_{B \cap E}|J T| d \lambda \leq t^{2 n} \sum_{i=1}^{\infty} \lambda\left[T\left(F_{i, j}\right)\right] \tag{11.36}
\end{equation*}
$$

Now let $j \rightarrow \infty$ and recall (11.33):

$$
\begin{aligned}
t^{-2 n} \int_{\mathbb{R}^{n}} N(T, B \cap E, y) d \lambda(y) & \leq \int_{B \cap E}|J T| d \lambda \\
& \leq t^{2 n} \int_{\mathbb{R}^{n}} N(T, B \cap E, y) d \lambda(y)
\end{aligned}
$$

Since $t$ was initially chosen as any number larger than 1 , we conclude that

$$
\begin{equation*}
\int_{B \cap E}|J T| d \lambda=\int_{\mathbb{R}^{n}} N(T, B \cap E, y) d \lambda(y) \tag{11.37}
\end{equation*}
$$

Now $B$ was defined as an arbitrary set of the sequence $\left\{B_{k}\right\}$ that was provided by Theorem 11.5. Thus (11.37) holds with $B$ replaced by an arbitrary $B_{k}$. Since the sets $\left\{B_{k}\right\}$ are disjoint and both sides of (11.37) are additive relative to $\cup B_{k}$, the Monotone Convergence Theorem implies

$$
\int_{E}|J T| d \lambda=\int_{\mathbb{R}^{n}} N(T, E, y) d \lambda(y)
$$

We have reached this conclusion under the assumption that $|J T| \neq 0$ on $E$. The proof will be concluded if we can show that

$$
\lambda\left[T\left(E_{0}\right)\right]=0
$$

where $E_{0}=E \cap\{x: J T(x)=0\}$. Note that $E_{0}$ is a Borel set. Note also that nothing is lost if we assume $E_{0}$ is bounded. Choose $\varepsilon>0$ and let $U \supset E_{0}$ be a bounded, open set with $\lambda\left(U-E_{0}\right)<\varepsilon$. For each $x \in E_{0}$, there exists $\delta_{x}>0$ such that $B\left(x, \delta_{x}\right) \subset U$ and

$$
|T(x+h)-T(x)|<\varepsilon|h|
$$

for all $h \in \mathbb{R}^{n}$ with $|h|<\delta_{x}$. In other words,

$$
\begin{equation*}
T[B(x, r)] \subset B(T(x), \varepsilon r) \tag{11.38}
\end{equation*}
$$

whenever $r<\delta_{x}$. The collection of all balls $B(x, r)$ where $x \in E_{0}$ and $r<\delta_{x}$ provides a Vitali covering of $E_{0}$ in the sense of Definition 7.4. Thus, by Theorem 7.7 , there exists a disjoint, countable subcollection $\left\{B_{i}\right\}$ such that

$$
\lambda\left(E_{0}-\bigcup_{i=1}^{\infty} B_{i}\right)=0
$$

Note that $\cup B_{i} \subset U$. Let us suppose that each $B_{i}$ is of the form $B_{i}=B\left(x_{i}, r_{i}\right)$. Then, using the fact that $T$ carries sets of measure zero into sets of measure zero,

$$
\begin{align*}
\lambda\left[T\left(E_{0}\right)\right] \leq & \sum_{i=1}^{\infty} \lambda\left[T\left(B_{i}\right)\right] \\
& \leq \sum_{i=1}^{\infty} \lambda\left[B\left(T\left(x_{i}\right), \varepsilon r_{i}\right)\right] \quad \text { by (11.38) }  \tag{11.38}\\
& =\sum_{i=1}^{\infty} \alpha(n)\left(\varepsilon r_{i}\right)^{n} \\
& =\varepsilon^{n} \sum_{i=1}^{\infty} \alpha(n) r_{i}^{n} \\
& =\varepsilon^{n} \sum_{i=1}^{\infty} \lambda\left(B_{i}\right) \\
& \leq \varepsilon^{n} \lambda(U) . \quad \text { since the }\left\{B_{i}\right\} \text { are disjoint }
\end{align*}
$$

Since $\lambda(U)<\infty$ (because $U$ is bounded) and $\varepsilon$ is arbitrary, we have $\lambda\left[T\left(E_{0}\right)\right]=$ 0.
11.8. Corollary. If $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is Lipschitz and univalent, then

$$
\int_{E}|J T(x)| d \lambda(x)=\lambda[T(E)]
$$

whenever $E \subset \mathbb{R}^{n}$ is measurable.
Proof. Since $T$ is univalent we have $N(T, E, y)=1$ on $T(E)$ and $N(T, E, y)=$ 0 in $\mathbb{R}^{n} \backslash T(E)$. Therefore, the result is clear from Theorem 11.7. See also Theorem 11.3.
11.9. ThEOREM. Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a Lipschitz transformation and suppose $f \in L^{1}\left(\mathbb{R}^{n}\right)$. Then

$$
\int_{\mathbb{R}^{n}} f \circ T(x)|J T(x)| d \lambda(x)=\int_{\mathbb{R}^{n}} f(y) N\left(T, \mathbb{R}^{n}, y\right) d \lambda(y)
$$

Proof. First, suppose $f$ is the characteristic function of a measurable set $A \subset \mathbb{R}^{n}$. Then $f \circ T=\chi_{T^{-1}(A)}$ and

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} f \circ T|J T| d \lambda=\int_{T^{-1}(A)}|J T| d \lambda & =\int_{\mathbb{R}^{n}} N\left(T, T^{-1}(A), y\right) d \lambda(y) \\
& =\int_{\mathbb{R}^{n}} \chi_{A}(y) N\left(T, \mathbb{R}^{n}, y\right) d \lambda(y) \\
& =\int_{\mathbb{R}^{n}} f(y) N\left(T, \mathbb{R}^{n}, y\right) d \lambda(y)
\end{aligned}
$$

Clearly, this also holds whenever $f$ is a simple function. Reference to Theorem 5.27 and the Monotone Convergence Theorem shows that it then holds whenever $f$ is nonnegative. Finally, writing $f=f^{+}-f^{-}$yields the final result.
11.10. Corollary. If $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is Lipschitz and univalent, then

$$
\int_{\mathbb{R}^{n}} f \circ T(x)|J T(x)| d \lambda(x)=\int_{\mathbb{R}^{n}} f(y) d \lambda(y)
$$

Proof. This is immediate from Theorem 11.9 since in this case $N\left(T, \mathbb{R}^{n}, y\right) \equiv$ 1.
11.11. Corollary. (Change of Variables Formula) Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a Lipschitz map and $f \in L^{1}\left(\mathbb{R}^{n}\right)$. If $E \subset \mathbb{R}^{n}$ is Lebesgue measurable and $T$ is injective on $E$, then

$$
\int_{T(E)} f(y) d \lambda(y)=\int_{E} f \circ T(x)|J T(x)| d \lambda(x)
$$

Proof. We apply Theorem 11.9 with $\chi_{T(E)} f$ instead of $f$.
11.12. Remark. This result provides a geometric interpretation of $J T(x)$. In case $T$ is linear, Theorem 11.3 states that $|J T|$ is given by

$$
\begin{equation*}
\frac{\lambda[T(E)]}{\lambda(E)} \tag{11.39}
\end{equation*}
$$

for any measurable set $E$. Roughly speaking, the same is true locally when $T$ is Lipschitz and univalent because Corollary 11.8 implies

$$
\int_{B\left(x_{0}, r\right)}|J T(x)| d \lambda(x)=\lambda\left[T\left(B\left(x_{0}, r\right)\right)\right]
$$

for an arbitrary ball $B\left(x_{0}, r\right)$. Now use the result on Lebesgue points (Theorem 7.11) to conclude that

$$
\left|J T\left(x_{0}\right)\right|=\lim _{r \rightarrow 0} \frac{\lambda\left[T\left(B\left(x_{0}, r\right)\right)\right]}{\lambda\left(B\left(x_{0}, r\right)\right)}
$$

for $\lambda$-almost all $x_{0}$, which is the infinitesimal analogue of (11.39).

We close this section with a discussion of spherical coordinates in $\mathbb{R}^{n}$. First, consider $\mathbb{R}^{2}$ with points designated by $\left(x_{1}, x_{2}\right)$. Let $\Omega$ denote $\mathbb{R}^{2}$ with the set $N:=\left\{\left(x_{1}, 0\right): x_{1} \geq 0\right\}$ removed. Let $r=\sqrt{x_{1}^{2}+x_{2}^{2}}$ and let $\theta$ be the angle from $N$ to the ray emanating from $(0,0)$ passing through $\left(x_{1}, x_{2}\right)$ with $0<\theta\left(x_{1}, x_{2}\right)<2 \pi$. Then $(r, \theta)$ are the coordinates of $\left(x_{1}, x_{2}\right)$ and

$$
x_{1}=r \cos \theta:=T^{1}(r, \theta), \quad x_{2}=r \sin \theta:=T^{2}(r, \theta)
$$

The transformation $T=\left(T^{1}, T^{2}\right)$ is a $C^{\infty}$ bijection of $(0, \infty) \times(0,2 \pi)$ onto $\mathbb{R}^{2}-N$. Furthermore, since $J T(r, \theta)=r$, we obtain from Corollary 11.11 that

$$
\int_{A} f \circ T(r, \theta) r d r d \theta=\int_{B} f d \lambda
$$

where $T(A)=B-N$. Of course, $\lambda(N)=0$ and therefore the integral over $B$ is the same as the integral over $T(A)$.

Now, we consider the case $n>2$ and proceed inductively. For $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$, Let $r=|x|$ and let $\theta_{1}=\cos ^{-1}\left(x_{1} / r\right), 0<\theta_{1}<\pi$. Also, let $\left(\rho, \theta_{2}, \ldots, \theta_{n-1}\right)$ be spherical coordinates for $x^{\prime}=\left(x_{2}, \ldots, x_{n}\right)$ where $\rho=\left|x^{\prime}\right|=r \sin \theta_{1}$. The coordinates of $x$ are

$$
\begin{aligned}
x_{1} & =r \cos \theta_{1}:=T^{1}\left(r, \theta_{1}, \ldots, \theta_{n}\right) \\
x_{2} & =r \sin \theta_{1} \cos \theta_{2}:=T^{2}\left(r, \theta_{1}, \ldots, \theta_{n}\right) \\
\vdots & \\
x_{n-1} & =r \sin \theta_{1} \sin \theta_{2} \ldots \sin \theta_{n-2} \cos \theta_{n-1}:=T^{n-1}\left(r, \theta_{1}, \ldots, \theta_{n}\right) \\
x_{n} & =r \sin \theta_{1} \sin \theta_{2} \ldots \sin \theta_{n-2} \sin \theta_{n-1}:=T^{n}\left(r, \theta_{1}, \ldots, \theta_{n}\right) .
\end{aligned}
$$

The mapping $T=\left(T^{1}, T^{2}, \ldots, T^{n}\right)$ is a bijection of

$$
(0, \infty) \times(0, \pi)^{n-2} \times(0,2 \pi)
$$

onto

$$
\mathbb{R}^{n}-\left(\mathbb{R}^{n-2} \times[0, \infty) \times\{0\}\right)
$$

A straightforward calculation shows that the Jacobian is

$$
J T\left(r, \theta_{1}, \ldots, \theta_{n-1}\right)=r^{n-1} \sin ^{n-2} \theta_{1} \sin ^{n-3} \theta_{2} \ldots \sin \theta_{n-2}
$$

Hence, with $\theta=\left(\theta_{1}, \ldots, \theta_{n-1}\right)$, we obtain

$$
\begin{align*}
\int_{A} f & \circ T(r, \theta) r^{n-1} \sin ^{n-2} \theta_{1} \sin ^{n-3} \theta_{2} \ldots \sin \theta_{n-2} d r d \theta_{1} \ldots \theta_{n-1} \\
& =\int_{T(A)} f d \lambda \tag{11.40}
\end{align*}
$$

## Exercises for Section 11.2

1. Prove that the sets $E(c, R, i)$ defined by (11.25) and (11.26) are Borel sets.
2. Give another proof of Theorem 11.4.

### 11.3. Sobolev Functions

It is not obvious that there is a class of functions defined on $\mathbb{R}^{n}$ that is analogous to the absolutely continuous functions on $\mathbb{R}$. However, it is tempting to employ the multi-dimensional analogue of Theorem 10.9 which states roughly that a function is absolutely continuous if and only if its derivative in the sense of distributions is a function. We will follow this direction and by considering functions whose partial derivatives are functions and show that this definition leads to a fruitful development.
11.13. Definition. Let $\Omega \subset \mathbb{R}^{n}$ be an open set and let $f \in L_{\text {loc }}^{1}(\Omega)$. We use Definition 10.7 to define the partial derivatives of $f$ in the sense of distributions. Thus, for $1 \leq i \leq n$, we say that a function $g_{i} \in L_{\text {loc }}^{1}(\Omega)$ is the $i^{\text {th }}$ partial derivative of $f$ in the sense of distributions if

$$
\begin{equation*}
\int_{\Omega} f \frac{\partial \varphi}{\partial x_{i}} d \lambda=-\int_{\Omega} g_{i} \varphi d \lambda \tag{11.41}
\end{equation*}
$$

for every test function $\varphi \in C_{c}^{\infty}(\Omega)$. We will write

$$
\frac{\partial f}{\partial x_{i}}=g_{i} ;
$$

thus, $\frac{\partial f}{\partial x_{i}}$ is merely defined to be the function $g_{i}$.
At this time we cannot assume that the partial derivative $\frac{\partial f}{\partial x_{i}}$ exists in the classical sense for the Sobolev function $f$. This existence will be discussed later after Theorem 11.19 has been proved. Consistent with the notation introduced in (10.5), we will sometimes write $D_{i} f$ to denote $\frac{\partial f}{\partial x_{i}}$.

The definition in (11.41) is a restatement of Definition 10.7 with the requirement that the derivative of $f$ (in the sense of distributions) is again a function and not merely a distribution. This requirement imposes a condition on $f$ and the purpose of this section is to see what properties $f$ must possess in order to satisfy this condition.

In general, we recall what it means for a higher order derivative of $f$ to be a function (see Definition 10.7). If $\alpha$ is a multi-index, then $g_{\alpha} \in L_{\text {loc }}^{1}(\Omega)$ is called the $\alpha^{\text {th }}$ distributional derivative of $f$ if

$$
\int_{\Omega} f D^{\alpha} \varphi d \lambda=(-1)^{|\alpha|} \int_{\Omega} \varphi g_{\alpha} d \lambda
$$

for every test function $\varphi \in C_{c}^{\infty}(\Omega)$. We write $D^{\alpha} f:=g_{\alpha}$. For $1 \leq p \leq \infty$ and $k$ a nonnegative integer, we say that $f$ belongs to the Sobolev Space

$$
W^{k, p}(\Omega)
$$

if $D^{\alpha} f \in L^{p}(\Omega)$ for each multi-index $\alpha$ with $|\alpha| \leq k$. In particular, this implies that $f \in L^{p}(\Omega)$. Similarly, the space

$$
W_{\mathrm{loc}}^{k, p}(\Omega)
$$

consists of all $f$ with $D^{\alpha} f \in L_{\text {loc }}^{p}(\Omega)$ for $|\alpha| \leq k$.
In order to motivate the definition of the distributional partial derivative, we recall the classical Gauss-Green Theorem.
11.14. Theorem. Suppose that $U$ is an open set with smooth boundary and let $\nu(x)$ denote the unit exterior normal to $U$ at $x \in \partial U$. If $V: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a transformation of class $C^{1}$, then

$$
\begin{equation*}
\int_{U} \operatorname{div} V d \lambda=\int_{\partial U} V(x) \cdot \nu(x) d H^{n-1}(x) \tag{11.42}
\end{equation*}
$$

where $\operatorname{div} V$, the divergence of $V=\left(V^{1}, \ldots, V^{n}\right)$, is defined by

$$
\operatorname{div} V=\sum_{i=1}^{n} \frac{\partial V^{i}}{\partial x_{i}}
$$

Now suppose that $f: \Omega \rightarrow \mathbb{R}$ is of class $C^{1}$ and $\varphi \in C_{c}^{\infty}(\Omega)$. Define $V: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ to be the transformation whose coordinate functions are all 0 except for the $i^{\text {th }}$ one, which is $f \varphi$. Then

$$
\operatorname{div} V=\frac{\partial(f \varphi)}{\partial x_{i}}=f \frac{\partial \varphi}{\partial x_{i}}+\frac{\partial f}{\partial x_{i}} \varphi
$$

Since the support of $\varphi$ is a compact set contained in $\Omega$, it is possible to find an open set $U$ with smooth boundary containing the support of $\varphi$ such that $\bar{U} \subset \Omega$. Then,

$$
\int_{\Omega} \operatorname{div} V d \lambda=\int_{U} \operatorname{div} V d \lambda=\int_{\partial U} V \cdot \nu d H^{n-1}=0
$$

Thus,

$$
\begin{equation*}
\int_{\Omega} f \frac{\partial \varphi}{\partial x_{i}} d \lambda=-\int_{\Omega} \frac{\partial f}{\partial x_{i}} \varphi d \lambda \tag{11.43}
\end{equation*}
$$

which is precisely (11.41) in case $f \in C^{1}(\Omega)$. Note that for a Sobolev function $f$, the formula (11.43) is valid for all test functions $\varphi \in C_{c}^{\infty}(\Omega)$ by the definition of the distributional derivative. The Sobolev norm of $f \in W^{1, p}(\Omega)$ is defined by

$$
\begin{equation*}
\|f\|_{1, p ; \Omega}:=\|f\|_{p ; \Omega}+\sum_{i=1}^{n}\left\|D_{i} f\right\|_{p ; \Omega} \tag{11.44}
\end{equation*}
$$

for $1 \leq p<\infty$ and

$$
\|f\|_{1, \infty ; \Omega}:=\operatorname{ess} \sup _{\Omega}\left(|f|+\sum_{i=1}^{n}\left|D_{i} f\right|\right)
$$

One can readily verify that $W^{1, p}(\Omega)$ becomes a Banach space when it is endowed with the above norm (Exercise 1, Section 11.3).
11.15. Remark. When $1 \leq p<\infty$, it can be shown (Exercise 2, Section 11.3) that the norm

$$
\begin{equation*}
\|f\|:=\|f\|_{p ; \Omega}+\left(\sum_{i=1}^{n}\left\|D_{i} f\right\|_{p ; \Omega}^{p}\right)^{1 / p} \tag{11.45}
\end{equation*}
$$

is equivalent to (11.44). Also, it is sometimes convenient to regard $W^{1, p}(\Omega)$ as a subspace of the Cartesian product of spaces $L^{p}(\Omega)$. The identification is made by defining $P: W^{1, p}(\Omega) \rightarrow \prod_{i=1}^{n+1} L^{p}(\Omega)$ as

$$
P(f)=\left(f, D_{1} f, \ldots, D_{n} f\right) \quad \text { for } \quad f \in W^{1, p}(\Omega)
$$

In view of Exercise 2, Section 8.1, it follows that $P$ is an isometric isomorphism of $W^{1, p}(\Omega)$ onto a subspace $W$ of this Cartesian product. Since $W^{1, p}(\Omega)$ is a Banach space, $W$ is a closed subspace. By Exercise 10, Section 8.4, $W$ is reflexive and therefore so is $W^{1, p}(\Omega)$.
11.16. Remark. Observe $f \in W^{1, p}(\Omega)$ is determined only up to a set of Lebesgue measure zero. We agree to call the Sobolev function $f$ continuous, bounded, etc. if there is a function $\bar{f}$ with these respective properties such that $\bar{f}=f$ almost everywhere.

We will show that elements in $W^{1, p}(\Omega)$ have representatives that permit us to regard them as generalizations of absolutely continuous functions on $\mathbb{R}^{1}$. First, we prove an important result concerning the convergence of regularizers of Sobolev functions.
11.17. Notation. If $\Omega^{\prime}$ and $\Omega$ are open sets of $\mathbb{R}^{n}$, we will write $\Omega^{\prime} \subset \subset \Omega$ to signify that the closure of $\Omega^{\prime}$ is compact and that the closure of $\Omega^{\prime}$ is a subset of $\Omega$. Also, we will frequently use $d x$ instead of $d \lambda(x)$ to denote integration with respect to Lebesgue measure.
11.18. Lemma. Suppose $f \in W^{1, p}(\Omega), 1 \leq p<\infty$. Then the mollifiers, $f_{\varepsilon}$, of $f$ (see (10.9)) satisfy

$$
\lim _{\varepsilon \rightarrow 0}\left\|f_{\varepsilon}-f\right\|_{1, p ; \Omega^{\prime}}=0
$$

whenever $\Omega^{\prime} \subset \subset \Omega$.

Proof. Since $\Omega^{\prime}$ is a bounded domain, there exists $\varepsilon_{0}>0$ such that $\varepsilon_{0}<$ $\operatorname{dist}\left(\Omega^{\prime}, \partial \Omega\right)$. For $\varepsilon<\varepsilon_{0}$, we may differentiate under the integral sign (see the proof
of (10.10)) to obtain for $x \in \Omega^{\prime}$ and $1 \leq i \leq n$,

$$
\begin{align*}
\frac{\partial f_{\varepsilon}}{\partial x_{i}}(x) & =\varepsilon^{-n} \int_{\Omega} \frac{\partial \varphi}{\partial x_{i}}\left(\frac{x-y}{\varepsilon}\right) f(y) d \lambda(y) \\
& =-\varepsilon^{-n} \int_{\Omega} \frac{\partial \varphi}{\partial y_{i}}\left(\frac{x-y}{\varepsilon}\right) f(y) d \lambda(y) \\
& =\varepsilon^{-n} \int_{\Omega} \varphi\left(\frac{x-y}{\varepsilon}\right) \frac{\partial f}{\partial y_{i}} d \lambda(y)  \tag{11.41}\\
& =\left(\frac{\partial f}{\partial x_{i}}\right)_{\varepsilon}(x)
\end{align*}
$$

Our result now follows from Theorem 10.1.

Since the definition of Sobolev functions requires that their distributional derivatives belong to $L^{p}$, it is natural to inquire if they have partial derivatives in the classical sense. To this end, we begin by showing that their partial derivatives exist almost everywhere. That is, in keeping with Remark 11.16, we will show that there is a function $f^{*}$ such that $f^{*}=f$ a.e. and that the partial derivatives of $f^{*}$ exist almost everywhere.

In the next theorem, the set $R$ is an interval of the form

$$
\begin{equation*}
R:=\left(a_{1}, b_{1}\right) \times \cdots \times\left(a_{n}, b_{n}\right) \tag{11.46}
\end{equation*}
$$

11.19. Theorem. Suppose $f \in W^{1, p}(\Omega), 1 \leq p<\infty$. Let $R \subset \subset \Omega$. Then $f$ has a representative $f^{*}$ that is absolutely continuous on almost all line segments of $R$ that are parallel to the coordinate axes, and the classical partial derivatives of $f^{*}$ agree almost everywhere with the distributional derivatives of $f$. Conversely, if $f=f^{*}$ almost everywhere for some function $f^{*}$ that is absolutely continuous on almost all line segments of $R$ that are parallel to the coordinate axes and if all partial derivatives of $f^{*}$ belong to $L^{p}(R)$, then $f \in W^{1, p}(R)$.

Proof. First, suppose $f \in W^{1, p}(\Omega)$ and fix $i$ with $1 \leq i \leq n$. Since $R \subset \subset \Omega$, the mollifiers $f_{\varepsilon}$ are defined for all $x \in R$ provided $\varepsilon$ is sufficiently small. Throughout the proof, only such mollifiers will be considered. We know from Lemma 11.18 that $\left\|f_{\varepsilon}-f\right\|_{1, p ; R} \rightarrow 0$ as $\varepsilon \rightarrow 0$. Choose a sequence $\varepsilon_{k} \rightarrow 0$ and let $f_{k}:=f_{\varepsilon_{k}}$. Also write $x \in R$ as $x=\left(x^{\prime}, t\right)$ where

$$
x^{\prime} \in R_{i}:=\left(a_{1}, b_{1}\right) \times \cdots \times \underset{\substack{\text { omitted }}}{\left(a_{i}, b_{i}\right) \times \cdots \times\left(a_{n}, b_{n}\right)}
$$

and $t \in\left(a_{i}, b_{i}\right), 1 \leq i \leq n$. Then it follows that

$$
\lim _{k \rightarrow \infty} \int_{R_{i}} \int_{a_{i}}^{b_{i}}\left|f_{k}\left(x^{\prime}, t\right)-f\left(x^{\prime}, t\right)\right|^{p}+\left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right|^{p} d t d \lambda\left(x^{\prime}\right)=0
$$

which can be rewriten as

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{R_{i}} F_{k}\left(x^{\prime}\right) d \lambda\left(x^{\prime}\right)=0 \tag{11.47}
\end{equation*}
$$

where

$$
F_{k}\left(x^{\prime}\right)=\int_{a_{i}}^{b_{i}}\left|f_{k}\left(x^{\prime}, t\right)-f\left(x^{\prime}, t\right)\right|^{p}+\left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right|^{p} d t
$$

As in the classical setting, we use $\nabla f$ to denote $\left(D_{1} f, \ldots, D_{n} f\right)$ which is defined in terms of the distributional derivatives of $f$. By Vitali's Convergence Theorem (see Theorem 6.30), and since (11.47) means $F_{k} \rightarrow 0$ in $L^{1}\left(R_{i}\right)$, there exists a subsequence (which will still be denoted as the full sequence) such that $F_{k}\left(x^{\prime}\right) \rightarrow 0$ for $\lambda_{n-1}$-a.e. $x^{\prime}$. That is,

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{a_{i}}^{b_{i}}\left|f_{k}\left(x^{\prime}, t\right)-f\left(x^{\prime}, t\right)\right|^{p}+\left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right|^{p} d t=0 \tag{11.48}
\end{equation*}
$$

for $\lambda_{n-1}$-almost all $x^{\prime} \in R_{i}$. From Hölder's inequality we have

$$
\begin{align*}
\int_{a_{i}}^{b_{i}} & \left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right| d t  \tag{11.49}\\
& \leq\left(b_{i}-a_{i}\right)^{1 / p^{\prime}}\left(\int_{a_{i}}^{b_{i}}\left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right|^{p} d t\right)^{1 / p}
\end{align*}
$$

The Fundamental Theorem of Calculus implies, for all $[a, b] \subset\left[a_{i}, b_{i}\right]$,

$$
\begin{aligned}
\left|f_{k}\left(x^{\prime}, b\right)-f_{k}\left(x^{\prime}, a\right)\right| & =\left|\int_{a}^{b} \frac{\partial f_{k}}{\partial x_{i}}\left(x^{\prime}, t\right) d t\right| \\
& \leq \int_{a}^{b}\left|\frac{\partial f_{k}}{\partial x_{i}}\left(x^{\prime}, t\right)\right| d t \\
& \leq \int_{a}^{b}\left|\nabla f_{k}\left(x^{\prime}, t\right)\right| d t \\
& \leq \int_{a}^{b}\left|\nabla f_{k}\left(x^{\prime}, t\right)-\nabla f\left(x^{\prime}, t\right)\right| d t+\int_{a}^{b}\left|\nabla f\left(x^{\prime}, t\right)\right| d t
\end{aligned}
$$

Consequently, it follows from (11.48), (11.49) and (11.50) that there is a constant $M_{x^{\prime}}$ such that

$$
\begin{equation*}
\left|f_{k}\left(x^{\prime}, b\right)-f_{k}\left(x^{\prime}, a\right)\right| \leq \int_{a_{i}}^{b_{i}}\left|\nabla f\left(x^{\prime}, t\right)\right| d t+1 \tag{11.51}
\end{equation*}
$$

whenever $k>M_{x^{\prime}}$ and $a, b \in\left[a_{i}, b_{i}\right]$. Note that (11.48) implies that there exists a subsequence of $\left\{f_{k}\left(x^{\prime}, \cdot\right)\right\}$, which again is denoted as the full sequence, such that

$$
\begin{equation*}
f_{k}\left(x^{\prime}, t\right) \rightarrow f\left(x^{\prime}, t\right), \quad \lambda \text {-a.e. t. } \tag{11.52}
\end{equation*}
$$

Fix $t_{0}$ for which $f_{k}\left(x^{\prime}, t_{0}\right) \rightarrow f_{k}\left(x^{\prime}, t_{0}\right), t_{0}<b$. Then (using this further subsequence), (11.51) with $a=t_{0}$ yields,

$$
\begin{equation*}
\left|f_{k}\left(x^{\prime}, b\right)\right| \leq 1+\int_{a_{i}}^{b_{i}}\left|\nabla f\left(x^{\prime}, t\right)\right| d t+\left|f_{k}\left(x^{\prime}, t_{0}\right)\right| \tag{11.53}
\end{equation*}
$$

Thus, for $k$ large enough (depending on $x^{\prime}$ ),

$$
\begin{equation*}
\left|f_{k}\left(x^{\prime}, b\right)\right| \leq 1+\int_{a_{i}}^{b_{i}}\left|\nabla f\left(x^{\prime}, t\right)\right| d t+\left|f\left(x^{\prime}, t_{0}\right)\right|+1 \tag{11.54}
\end{equation*}
$$

which proves that the sequence of functions $\left\{f_{k}\left(x^{\prime}, \cdot\right\}\right.$ is pointwise bounded for almost every $x^{\prime}$.

We now note that, for each $x^{\prime}$ under consideration, the functions $f_{k}\left(x^{\prime}, \cdot\right)$, as functions of $t$, are absolutely continuous (see Theorem 7.36). Moreover, the functions are absolutely continuous uniformly in $k$. Indeed, it is easy to check, using (11.50), that for each $\varepsilon>0$ there exists $\delta>0$ such that for any finite collection $\mathcal{F}$ of nonoverlapping intervals in $\left[a_{i}, b_{i}\right]$ with $\sum_{I \in \mathcal{F}}\left|b_{I}-a_{I}\right|<\delta$,

$$
\sum_{I \in \mathcal{F}}\left|f_{k}\left(x^{\prime}, b_{I}\right)-f_{k}\left(x^{\prime}, a_{I}\right)\right|<\varepsilon
$$

for all positive integers $k$. Here, as in Definition 7.23 , the endpoints of the interval $I$ are denoted by $a_{I}, b_{I}$. In particular, the sequence is equicontinuous on $\left[a_{i}, b_{i}\right]$. Since the sequence $\left\{f_{k}\left(x^{\prime}, \cdot\right)\right\}$ is pointwise bounded and equicontinuous, we now use the Arzelà-Ascoli Theorem, and we find that there is a subsequence that converges uniformly to a function on $\left[a_{i}, b_{i}\right]$, call it $f_{i}^{*}\left(x^{\prime}, \cdot\right)$. The uniform absolute continuity of this subsequence implies that $f_{i}^{*}\left(x^{\prime}, \cdot\right)$ is absolutely continuous. We now recall (11.52) and conclude that

$$
f\left(x^{\prime}, t\right)=f_{i}^{*}\left(x^{\prime}, t\right) \text { for } \lambda \text {-a.e. } t \in\left[a_{i}, b_{i}\right]
$$

To summarize what has been done so far, recall that for each interval $R \subset \Omega$ of the form (11.46) and each $1 \leq i \leq n$, there is a representative of $f, f_{i}^{*}$, that is absolutely continuous in $t$ for $\lambda_{n-1}$-almost all $x^{\prime} \in R_{i}$. This representative was obtained as the pointwise a.e. limit of a subsequence of mollifiers of $f$. Observe that there is a single subsequence and a single representative that can be found for all $i$ simultaneously. This can be accomplished by first obtaining a subsequence and a representative for $i=1$. Then a subsequence can be extracted from this one that can be used to define a representative for $i=1$ and 2 . Continuing in this way, the desired subsequence and representative are obtained. Thus, there is a sequence of mollifiers, denoted by $\left\{f_{k}\right\}$, and a function $f^{*}$ such that for each $1 \leq i \leq n, f_{k}\left(x^{\prime}, \cdot\right)$ converges uniformly to $f^{*}\left(x^{\prime}, \cdot\right)$ for $\lambda_{n-1}$-almost all $x^{\prime}$. Furthermore, $f^{*}\left(x^{\prime}, t\right)$ is absolutely continuous in $t$ for $\lambda_{n-1}$-almost all $x^{\prime}$.

The proof that classical partial derivatives of $f^{*}$ agree with the distributional derivatives almost everywhere is not difficult. Consider any partial derivative, say the $i^{\text {th }}$ one, and recall that $D_{i}=\frac{\partial}{\partial x_{i}}$. The distributional derivative is defined by

$$
D_{i} f(\varphi)=-\int_{R} f D_{i} \varphi d x
$$

for any test function $\varphi \in C_{c}^{\infty}(R)$. Since $f^{*}$ is absolutely continuous on almost all line segments parallel to the $i^{\text {th }}$ axis, we have

$$
\begin{aligned}
\int_{R} f^{*} D_{i} \varphi d x & =\int_{R_{i}} \int_{a_{i}}^{b_{i}} f^{*} D_{i} \varphi d x_{i} d x^{\prime} \\
& =-\int_{R_{i}} \int_{a_{i}}^{b_{i}}\left(D_{i} f^{*}\right) \varphi d x_{i} d x^{\prime} \\
& =-\int_{R}\left(D_{i} f^{*}\right) \varphi d x
\end{aligned}
$$

Since $f=f^{*}$ almost everywhere, we have

$$
\int_{R} f D_{i} \varphi d x=\int_{R} f^{*} D_{i} \varphi d x
$$

and therefore,

$$
\int_{R} D_{i} f \varphi d x=\int_{R} D_{i} f^{*} \varphi d x
$$

for every $\varphi \in C_{c}^{\infty}(R)$. This implies that the partial derivatives of $f^{*}$ agree with the distributional derivatives almost everywhere (see Exercise 3, Section 10.1). To prove the converse, suppose that $f$ has a representative $f^{*}$ as in the statement of the theorem. Then, for $\varphi \in C_{c}^{\infty}(R), f^{*} \varphi$ has the same absolutely continuous properties as does $f^{*}$. Thus, for $1 \leq i \leq n$, we can apply the Fundamental Theorem of Calculus to obtain

$$
\int_{a_{i}}^{b_{i}} \frac{\partial\left(f^{*} \varphi\right)}{\partial x_{i}}\left(x^{\prime}, t\right) d t=0
$$

for $\lambda_{n-1}$-almost all $x^{\prime} \in R_{i}$ and therefore (see problem 10, Section 7.5),

$$
\int_{a_{i}}^{b_{i}} f^{*}\left(x^{\prime}, t\right) \frac{\partial \varphi}{\partial x_{i}}\left(x^{\prime}, t\right) d t=-\int_{a_{i}}^{b_{i}} \frac{\partial f^{*}}{\partial x_{i}}\left(x^{\prime}, t\right) \varphi\left(x^{\prime}, t\right) d t
$$

Fubini's Theorem implies

$$
-\int_{R} f^{*} \frac{\partial \varphi}{\partial x_{i}} d \lambda=\int_{R} \frac{\partial f^{*}}{\partial x_{i}} \varphi d \lambda
$$

for all $\varphi \in C_{c}^{\infty}(R)$. Recall that the distributional derivative $\frac{\partial f}{\partial x_{i}}$ is defined as

$$
\frac{\partial f}{\partial x_{i}}(\varphi)=-\int_{R} f \frac{\partial \varphi}{\partial x_{i}} d \lambda, \quad \varphi \in C_{c}^{\infty}(R)
$$

Hence, since $f=f^{*}$ a.e.,

$$
\begin{equation*}
\frac{\partial f}{\partial x_{i}}(\varphi)=\int_{R} \frac{\partial f^{*}}{\partial x_{i}} \varphi d \lambda, \quad \text { for all } C_{c}^{\infty}(R) \tag{11.55}
\end{equation*}
$$

From (11.55), it follows that the distribution $\frac{\partial f}{\partial x_{i}}$ is the function $\frac{\partial f^{*}}{\partial x_{i}} \in L^{p}(R)$. Therefore $\frac{\partial f}{\partial x_{i}} \in L^{p}(R)$ and we conclude $f \in W^{1, p}(R)$.

## Exercises for Section 11.3

1. Prove that $W^{1, p}(\Omega)$ endowed with the norm (11.44) is a Banach space.
2. Prove that the norm

$$
\|f\|:=\|f\|_{p ; \Omega}+\left(\sum_{i=1}^{n}\left\|D_{i} f\right\|_{p ; \Omega}^{p}\right)^{1 / p}
$$

is equivalent to (11.44).
3. With the help of Exercise 2, Section 8.1, show that $W^{1, p}(\Omega)$ can be regarded as a closed subspace of the Cartesian product of $L^{p}(\Omega)$ spaces. Referring now to 8.10 , show that this subspace is reflexive and hence $W^{1, p}(\Omega)$ is reflexive if $1<p<\infty$.
4. Suppose $f \in W^{1, p}\left(\mathbb{R}^{n}\right)$. Prove that $f^{+}$and $f^{-}$are in $W^{1, p}\left(\mathbb{R}^{n}\right)$. Hint: use Theorem 11.19.

### 11.4. Approximating Sobolev Functions

> We will show that the Sobolev space $W^{1, p}(\Omega)$ can be characterized as the closure of $C^{\infty}(\Omega)$ in the Sobolev norm. This is a very useful result and is employed frequently in applications. In the next section we will demonstrate its utility in proving the Sobolev inequality, which implies that $W^{1, p}(\Omega) \subset L^{p *}(\Omega)$, where $p *=n p /(n-p)$.

We begin with a smooth version of Theorem 9.8.
11.20. Lemma. Let $\mathcal{G}$ be an open cover of a set $E \subset \mathbb{R}^{n}$. Then there exists a family $\mathcal{F}$ of functions $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $0 \leq f \leq 1$ and
(i) For each $f \in \mathcal{F}$, there exists $U \in \mathcal{G}$ such that $\operatorname{spt} f \subset U$,
(ii) If $K \subset E$ is compact, then spt $f \cap K \neq 0$ for only finitely many $f \in \mathcal{F}$,
(iii) $\sum_{f \in \mathcal{F}} f(x)=1$ for each $x \in E$. The family $\mathcal{F}$ is called a smooth partition of unity of $E$ subordinate to the open covering $\mathcal{G}$.

Proof. In case $E$ is compact, our desired result follows from the proof of Theorem 9.8 and Exercise 1, Section 10.1.

Now assume $E$ is open and for each positive integer $i$ define

$$
E_{i}=E \cap \bar{B}(0, i) \cap\left\{x: \operatorname{dist}(x, \partial E) \geq \frac{1}{i}\right\}
$$

Thus, $E_{i}$ is compact, $E_{i} \subset \operatorname{int} E_{i+1}$, and $E=\cup_{i=1}^{\infty} E_{i}$. Let $\mathcal{G}_{i}$ denote the collection of all open sets of the form

$$
U \cap\left\{\operatorname{int} E_{i+1}-E_{i-2}\right\}
$$

where $U \in \mathcal{G}$ and where we take $E_{0}=E_{-1}=\emptyset$. The family $\mathcal{G}_{i}$ forms an open covering of the compact set $E_{i}-\operatorname{int} E_{i-1}$. Therefore, our first case applies and we obtain a smooth partition of unity, $\mathcal{F}_{i}$, subordinate to $\mathcal{G}_{i}$, which consists of finitely many elements. For each $x \in \mathbb{R}^{n}$, define

$$
s(x)=\sum_{i=1}^{\infty} \sum_{g \in \mathcal{F}_{i}} g(x)
$$

The sum is well defined since only a finite number of terms are nonzero for each $x$. Note that $s(x)>0$ for $x \in E$. Corresponding to each positive integer $i$ and to each function $g \in \mathcal{G}_{i}$, define a function $f$ by

$$
f(x)= \begin{cases}\frac{g(x)}{s(x)} & \text { if } x \in E \\ 0 & \text { if } x \notin E\end{cases}
$$

The partition of unity $\mathcal{F}$ that we want is comprised of all such functions $f$.
If $E$ is arbitrary, then any partition of unity for the open set $\cup\{U: U \in \mathcal{G}\}$ is also one for $E$.

Clearly, the set

$$
S=C^{\infty}(\Omega) \cap\left\{f:\|f\|_{1, p ; \Omega}<\infty\right\}
$$

is contained in $W^{1, p}(\Omega)$. Moreover, the same is true of the closure of $S$ in the Sobolev norm since $W^{1, p}(\Omega)$ is complete. The next result shows that $\bar{S}=W^{1, p}(\Omega)$.
11.21. Theorem. If $1 \leq p<\infty$, then the space

$$
C^{\infty}(\Omega) \cap\left\{f:\|f\|_{1, p ; \Omega}<\infty\right\}
$$

is dense in $W^{1, p}(\Omega)$.
Proof. Let $\Omega_{i}$ be subdomains of $\Omega$ such that $\overline{\Omega_{i}} \subset \Omega_{i+1}$ and $\cup_{i=1}^{\infty} \Omega_{i}=\Omega$. Let $\mathcal{F}$ be a partition of unity of $\Omega$ subordinate to the covering $\left\{\Omega_{i+1}-\overline{\Omega_{i-1}}\right\}$, where $\Omega_{-1}$ is defined as the null set. Let $f_{i}$ denote the sum of the finitely many $f \in \mathcal{F}$ with spt $f \subset \Omega_{i+1}-\overline{\Omega_{i-1}}$. Then $f_{i} \in C_{c}^{\infty}\left(\Omega_{i+1}-\overline{\Omega_{i-1}}\right)$ and

$$
\begin{equation*}
\sum_{i=1} f_{i} \equiv 1 \quad \text { on } \quad \Omega \tag{11.56}
\end{equation*}
$$

Choose $\varepsilon>0$. For $f \in W^{1, p}(\Omega)$, refer to Lemma 11.18 to obtain $\varepsilon_{i}>0$ such that

$$
\begin{align*}
& \operatorname{spt}\left(f_{i} f\right)_{\varepsilon_{i}} \subset \Omega_{i+1}-\overline{\Omega_{i-1}}  \tag{11.57}\\
& \left\|\left(f_{i} f\right)_{\varepsilon_{i}}-f_{i} f\right\|_{1, p ; \Omega}<\varepsilon 2^{-i}
\end{align*}
$$

With $g_{i}:=\left(f_{i} f\right)_{\varepsilon_{i}}$, (11.57) implies that only finitely many of the $g_{i}$ can fail to vanish on any given $\Omega^{\prime} \subset \subset \Omega$. Therefore $g:=\sum_{i=1}^{\infty} g_{i}$ is defined and is an element of $C^{\infty}(\Omega)$. For $x \in \Omega_{i}$, we have

$$
f(x)=\sum_{j=1}^{i} f_{j}(x) f(x)
$$

and by (11.57)

$$
g(x)=\sum_{j=1}^{i}\left(f_{j} f\right)_{\varepsilon_{j}}(x)
$$

Consequently,

$$
\|f-g\|_{1, p ; \Omega_{i}} \leq \sum_{j=1}^{i}\left\|\left(f_{j} f\right)_{\varepsilon_{j}}-f_{j} f\right\|_{1, p ; \Omega}<\varepsilon
$$

Now an application of the Monotone Convergence Theorem establishes our desired result.

The previous result holds in particular when $\Omega=\mathbb{R}^{n}$, in which case we get the following apparently stronger result.
11.22. Corollary. If $1 \leq p<\infty$, the space $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ is dense in $W^{1, p}\left(\mathbb{R}^{n}\right)$.

Proof. This follows from the previous result and the fact that $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ is dense in

$$
C^{\infty}\left(\mathbb{R}^{n}\right) \cap\left\{f:\|f\|_{1, p}<\infty\right\}
$$

relative to the Sobolev norm (see Exercise 1, Section 11.4).
Recall that if $f \in L^{p}\left(\mathbb{R}^{n}\right)$, then $\|f(x+h)-f(x)\|_{p} \rightarrow 0$ as $h \rightarrow 0$. A similar result provides a very useful characterization of $W^{1, p}$.
11.23. Theorem. Let $1<p<\infty$ and suppose $\Omega \subset \mathbb{R}^{n}$ is an open set. If $f \in W^{1, p}(\Omega)$ and $\Omega^{\prime} \subset \subset \Omega$, then $\left|h^{-1}\right|\|f(x+h)-f(x)\|_{p ; \Omega^{\prime}}$ remains bounded for all sufficiently small $h$. Conversely, if $f \in L^{p}(\Omega)$ and

$$
\left|h^{-1}\right|\|f(x+h)-f(x)\|_{p ; \Omega^{\prime}}
$$

remains bounded for all sufficiently small $h$, then $f \in W^{1, p}\left(\Omega^{\prime}\right)$.
Proof. Assume $f \in W^{1, p}(\Omega)$ and let $\Omega^{\prime} \subset \subset \Omega$. By Theorem 11.21, there exists a sequence of $C^{\infty}(\Omega)$ functions $\left\{f_{k}\right\}$ such that $\left\|f_{k}-f\right\|_{1, p ; \Omega} \rightarrow 0$ as $k \rightarrow \infty$. For each $g \in C^{\infty}(\Omega)$, we have

$$
\frac{g(x+h)-g(x)}{|h|}=\frac{1}{|h|} \int_{0}^{|h|} \nabla g\left(x+t \frac{h}{|h|}\right) \cdot \frac{h}{|h|} d t
$$

so by Jensen's inequality (Exercise 10, Section 6.5),

$$
\left|\frac{g(x+h)-g(x)}{h}\right|^{p} \leq \frac{1}{|h|} \int_{0}^{|h|}\left|\nabla g\left(x+t \frac{h}{|h|}\right)\right|^{p} d t
$$

whenever $x \in \Omega^{\prime}$ and $h<\delta:=\operatorname{dist}\left(\partial \Omega^{\prime}, \partial \Omega\right)$. Therefore,

$$
\begin{aligned}
\|g(x+h)-g(x)\|_{p ; \Omega^{\prime}}^{p} & \leq|h|^{p} \frac{1}{|h|} \int_{0}^{|h|} \int_{\Omega^{\prime}}\left|\nabla g\left(x+t \frac{h}{|h|}\right)\right|^{p} d x d t \\
& \leq|h|^{p-1} \int_{0}^{|h|} \int_{\Omega}|\nabla g(x)|^{p} d x d t
\end{aligned}
$$

or

$$
\|g(x+h)-g(x)\|_{p ; \Omega^{\prime}} \leq|h|\|\nabla g\|_{p ; \Omega}
$$

for all $h<\delta$. As this inequality holds for each $f_{k}$, it also holds for $f$.
For the proof of the converse, let $e_{i}$ denote the $i^{\text {th }}$ unit basis vector. By assumption, the sequence

$$
\left\{\frac{f\left(x+e_{i} / k\right)-f(x)}{1 / k}\right\}
$$

is bounded in $L^{p}\left(\Omega^{\prime}\right)$ for all large $k$. Therefore by Alaoglu's Theorem (Theorem 8.40), there exist a subsequence (denoted by the full sequence) and $f_{i} \in L^{p}\left(\Omega^{\prime}\right)$ such that

$$
\frac{f\left(x+e_{i} / k\right)-f(x)}{1 / k} \rightarrow f_{i}
$$

weakly in $L^{p}\left(\Omega^{\prime}\right)$. Thus, for $\varphi \in C_{0}^{\infty}\left(\Omega^{\prime}\right)$,

$$
\begin{aligned}
\int_{\Omega^{\prime}} f_{i} \varphi d x & =\lim _{k \rightarrow \infty} \int_{\Omega^{\prime}}\left[\frac{f\left(x+e_{i} / k\right)-f(x)}{1 / k}\right] \varphi(x) d x \\
& =\lim _{k \rightarrow \infty} \int_{\Omega^{\prime}} f(x)\left[\frac{\varphi\left(x-e_{i} / k\right)-\varphi(x)}{1 / k}\right] d x \\
& =-\int_{\Omega^{\prime}} f D_{i} \varphi d x
\end{aligned}
$$

This shows that $D_{i} f=f_{i}$ in the sense of distributions. Hence, $f \in W^{1, p}\left(\Omega^{\prime}\right)$.

## Exercises for Section 11.4

1. Relative to the Sobolev norm, prove that $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ is dense in

$$
C^{\infty}\left(\mathbb{R}^{n}\right) \cap W^{1, p}\left(\mathbb{R}^{n}\right)
$$

2. Let $f \in C^{2}(\Omega)$. Show that $f$ is harmonic in $\Omega$ if and only if $f$ is weakly harmonic in $\Omega$ (i.e., $\int_{\Omega} f \Delta \varphi d x=0$ for every $\left.\varphi \in C_{c}^{\infty}(\Omega)\right)$. Show also that $f$ is weakly harmonic in $\Omega$ if and only if $\int_{\Omega} \nabla \varphi \cdot \nabla f d x=0$ for every $\varphi \in C_{c}^{\infty}(\Omega)$.
3. Let $f \in W^{1,2}(\Omega)$. Show that $f$ is weakly harmonic (i.e.,

$$
\int_{\Omega} f \Delta \varphi d x=0
$$

for every $\left.\varphi \in C_{c}^{\infty}(\Omega)\right)$ if and only if

$$
\int_{\Omega} \nabla f \cdot \nabla \varphi d x=0
$$

for every test function $\varphi$. Hint: Use Theorem 11.21 along with (11.42).
4. Let $\Omega \subset \mathbb{R}^{n}$ be an open connected set and suppose $f \in W^{1, p}(\Omega)$ has the property that $\nabla f=0$ almost everywhere in $\Omega$. Prove that $f$ is constant on $\Omega$.

### 11.5. Sobolev Imbedding Theorem

One of the most useful estimates in the theory of Sobolev functions is the Sobolev inequality. It implies that a Sobolev function is in a higher Lebesgue class than the one in which it was originally defined. In fact, for $1 \leq p<n, W^{1, p}(\Omega) \subset L^{p^{*}}(\Omega)$ where $p^{*}=n p /(n-p)$.

We use the result of the previous section to set the stage for the next definition. First, consider a bounded open set $\Omega \subset \mathbb{R}^{n}$ whose boundary has Lebesgue measure zero. Recall that Sobolev functions are only defined almost everywhere. Consequently, it is not possible in our present state of development to define what it means for a Sobolev function to be zero (pointwise) on the boundary of a domain $\Omega$. Instead, we define what it means for a Sobolev function to be zero on $\partial \Omega$ in a global sense.
11.24. Definition. Let $\Omega \subset \mathbb{R}^{n}$ be an arbitrary open set. The space $W_{0}^{1, p}(\Omega)$ is defined as the closure of $C_{c}^{\infty}(\Omega)$ relative to the Sobolev norm. Thus, $f \in W_{0}^{1, p}(\Omega)$ if and only if there is a sequence of functions $f_{k} \in C_{c}^{\infty}(\Omega)$ such that

$$
\lim _{k \rightarrow \infty}\left\|f_{k}-f\right\|_{1, p ; \Omega}=0
$$

11.25. Remark. If $\Omega=\mathbb{R}^{n}$, Exercise 1, Section 11.4, implies that $W_{0}^{1, p}\left(\mathbb{R}^{n}\right)=$ $W^{1, p}\left(\mathbb{R}^{n}\right)$.
11.26. Theorem. Let $1 \leq p<n$ and let $\Omega \subset \mathbb{R}^{n}$ be any open set. There is a constant $C=C(n, p)$ such that for $f \in W_{0}^{1, p}(\Omega)$,

$$
\|f\|_{p^{*} ; \Omega} \leq C\|\nabla f\|_{p ; \Omega}
$$

Proof. Step 1: Assume first that $p=1$ and $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. Appealing to the Fundamental Theorem of Calculus and using the fact that $f$ has compact support, it follows for each integer $i, 1 \leq i \leq n$, that

$$
f\left(x_{1}, \ldots, x_{i}, \ldots, x_{n}\right)=\int_{-\infty}^{x_{i}} \frac{\partial f}{\partial x_{i}}\left(x_{1}, \ldots, t_{i}, \ldots, x_{n}\right) d t_{i}
$$

and therefore,

$$
\begin{aligned}
|f(x)| & \leq \int_{-\infty}^{\infty}\left|\frac{\partial f}{\partial x_{i}}\left(x_{1}, \ldots, t_{i}, \ldots, x_{n}\right)\right| d t_{i} \\
& \leq \int_{-\infty}^{\infty}\left|\nabla f\left(x_{1}, \ldots, t_{i}, \ldots, x_{n}\right)\right| d t_{i}, \quad 1 \leq i \leq n
\end{aligned}
$$

Consequently,

$$
|f(x)|^{\frac{n}{n-1}} \leq \prod_{i=1}^{n}\left(\int_{-\infty}^{\infty}\left|\nabla f\left(x_{1}, \ldots, t_{i}, \ldots, x_{n}\right)\right| d t_{i}\right)^{\frac{1}{n-1}}
$$

This can be rewritten as

$$
\begin{aligned}
& |f(x)|^{\frac{n}{n-1}} \\
& \leq\left(\int_{-\infty}^{\infty}|\nabla f(x)| d t_{1}\right)^{\frac{1}{n-1}} \cdot \prod_{i=2}^{n}\left(\int_{-\infty}^{\infty}\left|\nabla f\left(x_{1}, \ldots, t_{i}, \ldots, x_{n}\right)\right| d t_{i}\right)^{\frac{1}{n-1}}
\end{aligned}
$$

Only the first factor on the right is independent of $x_{1}$. Thus, when the inequality is integrated with respect to $x_{1}$ we obtain, with the help of generalized Hölder's inequality (see Exercise 14, Section 6.5),

$$
\begin{aligned}
& \int_{-\infty}^{\infty}|f|^{\frac{n}{n-1}} d x_{1} \\
& \quad \leq\left(\int_{-\infty}^{\infty}|\nabla f| d t_{1}\right)^{\frac{1}{n-1}} \int_{-\infty}^{\infty} \prod_{i=2}^{n}\left(\int_{-\infty}^{\infty}|\nabla f| d t_{i}\right)^{\frac{1}{n-1}} d x_{1} \\
& \quad \leq\left(\int_{-\infty}^{\infty}|\nabla f| d t_{1}\right)^{\frac{1}{n-1}}\left(\prod_{i=2}^{n} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d x_{1} d t_{i}\right)^{\frac{1}{n-1}}
\end{aligned}
$$

Similarly, integration with respect to $x_{2}$ yields

$$
\begin{aligned}
& \int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|f|^{\frac{n}{n-1}} d x_{1} d x_{2} \\
& \quad \leq\left(\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d x_{1} d t_{2}\right)^{\frac{1}{n-1}} \int_{-\infty}^{\infty} \prod_{i=1, i \neq 2}^{n} I_{i}^{\frac{1}{n-1}} d x_{2} \\
& I_{1}:=\int_{-\infty}^{\infty}|\nabla f| d t_{1} \quad I_{i}=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d x_{1} d t_{i}, \quad i=3, \ldots, n
\end{aligned}
$$

Applying once more the generalized Hölder inequality, we find

$$
\begin{aligned}
& \int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|f|^{\frac{n}{n-1}} d x_{1} d x_{2} \\
& \quad \leq\left(\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d x_{1} d t_{2}\right)^{\frac{1}{n-1}}\left(\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d t_{1} d x_{2}\right)^{\frac{1}{n-1}} \\
& \quad \times \prod_{i=3}^{n}\left(\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty}|\nabla f| d x_{1} d x_{2} d t_{i}\right)^{\frac{1}{n-1}}
\end{aligned}
$$

Continuing this way for the remaining $n-2$ steps, we finally arrive at

$$
\int_{\mathbb{R}^{n}}|f|^{\frac{n}{n-1}} d x \leq \prod_{i=1}^{n}\left(\int_{\mathbb{R}^{n}}|\nabla f| d x\right)^{\frac{1}{n-1}}
$$

or

$$
\begin{equation*}
\|f\|_{\frac{n}{n-1}} \leq \int_{\mathbb{R}^{n}}|\nabla f| d x, \quad f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right) \tag{11.58}
\end{equation*}
$$

which is the desired result in the case $p=1$ and $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$.
Step 2: Assume now that $1 \leq p<n$ and $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. This case is treated by replacing $f$ by positive powers of $|f|$. Thus, for $q$ to be determined later, apply our previous step to $g:=|f|^{q}$. Technically, the previous step requires $g \in C_{c}^{\infty}(\Omega)$. However, a close examination of the proof reveals that we only need $g$ to be an absolutely continuous function in each variable separately. Then,

$$
\begin{aligned}
\left\|f^{q}\right\|_{n /(n-1)} & \leq\left.\int_{\mathbb{R}^{n}}|\nabla| f\right|^{q} \mid d x \\
& =\int_{\mathbb{R}^{n}} q|f|^{q-1}|\nabla f| d x \\
& \leq q\left\|f^{q-1}\right\|_{p^{\prime}}\|\nabla f\|_{p}
\end{aligned}
$$

where we have used Hölder's inequality in the last inequality. Requesting $(q-1) p^{\prime}=$ $q \frac{n}{n-1}$, with $\frac{1}{p^{\prime}}=1-\frac{1}{p}$, we obtain $q=(n-1) p /(n-p)$. With this $q$ we have $q \frac{n}{n-1}=\frac{n p}{n-p}$ and hence

$$
\left(\int_{\mathbb{R}^{n}}|f|^{\frac{n p}{n-p}}\right)^{\frac{n-1}{n}} \leq q\left(\int_{\mathbb{R}^{n}}|f|^{\frac{n p}{n-p}}\right)^{\frac{1}{p^{\prime}}}\|\nabla f\|_{p ; \mathbb{R}^{n}}
$$

Since $\frac{n-1}{n}-\frac{1}{p^{\prime}}=\frac{n-p}{n p}$ we obtain

$$
\left(\int_{\mathbb{R}^{n}}|f|^{\frac{n p}{n-p}}\right)^{\frac{n-p}{n p}} \leq q\|\nabla f\|_{p ; \mathbb{R}^{n}}
$$

which is

$$
\begin{equation*}
\|f\|_{\frac{n p}{n-p} ; \mathbb{R}^{n}} \leq \frac{(n-1) p}{n-p}\|\nabla f\|_{p ; \mathbb{R}^{n}}, f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right) \tag{11.59}
\end{equation*}
$$

Step 3:

Let $f \in W_{0}^{1, p}(\Omega)$. Let $\left\{f_{i}\right\}$ be a sequence of functions in $C_{c}^{\infty}(\Omega)$ converging to $f$ in the Sobolev norm. We have

$$
\begin{equation*}
\left\|f_{i}-f_{j}\right\|_{1, p ; \Omega} \rightarrow 0 \tag{11.60}
\end{equation*}
$$

Applying (11.59) to $\left\|f_{i}-f_{j}\right\|$ we obtain

$$
\left\|f_{i}-f_{j}\right\|_{p^{*} ; \mathbb{R}^{n}} \leq C\left\|f_{i}-f_{j}\right\|_{1, p ; \mathbb{R}^{n}},
$$

where we have extended the functions $f_{i}$ by zero outside $\Omega$. Therefore we have

$$
\begin{equation*}
\left\|f_{i}-f_{j}\right\|_{p^{*} ; \Omega} \leq C\left\|f_{i}-f_{j}\right\|_{1, p ; \Omega} . \tag{11.61}
\end{equation*}
$$

From (11.60) and (11.61) it follows that $\left\{f_{i}\right\}$ is Cauchy in $L^{p^{*}}(\Omega)$, and hence there exists $g \in L^{p^{*}}(\Omega)$ such that

$$
f_{i} \rightarrow g \text { in } L^{p^{*}}(\Omega) .
$$

Therefore, there exists a subsequence $f_{i_{k}}$ such that $f_{i_{k}} \rightarrow g$ pointwise. Since $f_{i_{k}} \rightarrow f$ in $L^{p}(\Omega)$ then, up to a further subsequence, $f_{i_{k}} \rightarrow f$ pointwise. By uniqueness of the limit, we conclude $f=g$ almost everywhere. That is,

$$
\begin{equation*}
f_{i} \rightarrow f \text { in } L^{p^{*}}(\Omega) . \tag{11.62}
\end{equation*}
$$

Using that $\left|\nabla f_{i}\right| \rightarrow|\nabla f|$ in $L^{p}(\Omega)$ and (11.62), we can let $i \rightarrow \infty$ in

$$
\left\|f_{i}\right\|_{p^{*} ; \Omega} \leq C\left\|\nabla f_{i}\right\|_{p ; \Omega}
$$

to obtain

$$
\|f\|_{p^{*} ; \Omega} \leq C\|\nabla f\|_{p ; \Omega} .
$$

## Exercises for Section 11.5

1. Suppose $\left\|f_{i}-f_{j}\right\|_{q ; \Omega} \rightarrow 0$ and $\left\|f_{i}-f\right\|_{p ; \Omega} \rightarrow 0$ where $\Omega \subset \mathbb{R}^{n}$ and $q>p$. Prove that $\left\|f_{i}-f\right\|_{q ; \Omega} \rightarrow 0$.

### 11.6. Applications

> A basic problem in the calculus of variations is to find a harmonic function in a domain that assumes values that have been prescribed on the boundary. Using results of the previous sections, we will discuss a solution to this problem. Our first result shows that this problem has a "weak solution," that is, a solution in the sense of distributions.
11.27. Definition. A function $f \in C^{2}(\Omega)$ is called harmonic in $\Omega$ if

$$
\frac{\partial^{2} f}{\partial x_{1}^{2}}(x)+\frac{\partial^{2} f}{\partial x_{2}^{2}}(x)+\cdots+\frac{\partial^{2} f}{\partial x_{n}^{2}}(x)=0
$$

for each $x \in \Omega$.

A straightforward calculation shows that this is equivalent to $\operatorname{div}(\nabla f)(x)=0$. In general, if we let $\Delta$ denote the operator

$$
\Delta=\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\cdots+\frac{\partial^{2}}{\partial x_{n}^{2}}
$$

taking $V=\varphi \nabla f$ in (11.42) we have

$$
\begin{equation*}
0=\int_{\Omega} \operatorname{div} V=\int_{\Omega} \varphi \Delta f d x+\int_{\Omega} \nabla \varphi \cdot \nabla f d x \tag{11.63}
\end{equation*}
$$

whenever $f \in C^{2}(\Omega)$ and $\varphi \in C_{c}^{\infty}(\Omega)$, and hence if $f$ is harmonic in $\Omega$ then

$$
\int_{\Omega} \nabla f \cdot \nabla \varphi d x=0
$$

whenever $\varphi \in C_{c}^{\infty}(\Omega)$.
Returning to the context of distributions, recall that a distribution can be differentiated any number of times. In particular, it is possible to define $\Delta T$ whenever $T$ is a distribution. If $T$ is taken as an integrable function $f$, we have

$$
\Delta f(\varphi)=\int_{\Omega} f \Delta \varphi d x
$$

for all $\varphi \in C_{c}^{\infty}(\Omega)$.
11.28. Definition. We say that $f$ is harmonic in the sense of distributions or weakly harmonic if $\Delta f(\varphi)=0$ for every test function $\varphi$.

Therefore (see Exercise 1, Section 11.4), $f \in C^{2}(\Omega)$ is harmonic in $\Omega$ if and only if $f$ is weakly harmonic in $\Omega$. Also, $f \in C^{2}(\Omega)$ is weakly harmonic in $\Omega$ if and only if $\int_{\Omega} \nabla \varphi \cdot \nabla f d x=0$.

If $f \in W^{1,2}(\Omega)$ is weakly harmonic, then (11.63) implies (with $f$ and $\varphi$ interchanged) that

$$
\begin{equation*}
\int_{\Omega} \nabla f \cdot \nabla \varphi d x=0 \tag{11.64}
\end{equation*}
$$

for every test function $\varphi \in C_{c}^{\infty}(\Omega)$ (see Exercise 2, Section 11.6). Since $f \in$ $W^{1,2}(\Omega)$, note that (11.64) remains valid with $\varphi \in W_{0}^{1,2}(\Omega)$.
11.29. Theorem. Suppose $\Omega \subset \mathbb{R}^{n}$ is a bounded open set and let $\psi \in W^{1,2}(\Omega)$. Then there exists a weakly harmonic function $f \in W^{1,2}(\Omega)$ such that $f-\psi \in$ $W_{0}^{1,2}(\Omega)$.

The theorem states that for a given function $\psi \in W^{1,2}(\Omega)$, there exists a weakly harmonic function $f$ that assumes the same values (in a weak sense) on $\partial \Omega$ as does $\psi$. That is, $f$ and $\psi$ have the same boundary values and

$$
\int_{\Omega} \nabla f \cdot \nabla \varphi d x=0
$$

for every test function $\varphi \in C_{c}^{\infty}(\Omega)$. Later we will show that $f$ is, in fact, harmonic in the sense of Definition 11.27. In particular, it will be shown that $f \in C^{\infty}(\Omega)$.

Proof. Let

$$
\begin{equation*}
m:=\inf \left\{\int_{\Omega}|\nabla f|^{2} d x: f-\psi \in W_{0}^{1,2}(\Omega)\right\} \tag{11.65}
\end{equation*}
$$

This definition requires $f-\psi \in W_{0}^{1,2}(\Omega)$. Since $\psi \in W^{1,2}(\Omega)$, note that $f$ must be an element of $W^{1,2}(\Omega)$ and therefore that $|\nabla f| \in L^{2}(\Omega)$.

Our first objective is to prove that the infimum is attained. For this purpose, let $f_{i}$ be a sequence such that $f_{i}-\psi \in W_{0}^{1,2}(\Omega)$ and

$$
\begin{equation*}
\int_{\Omega}\left|\nabla f_{i}\right|^{2} d x \rightarrow m \quad \text { as } \quad i \rightarrow \infty \tag{11.66}
\end{equation*}
$$

Now apply both Hölder's and Sobolev's inequality (Theorem 11.26) to obtain

$$
\left\|f_{i}-\psi\right\|_{2 ; \Omega} \leq \lambda(\Omega)^{\left(\frac{1}{2}-\frac{1}{2^{*}}\right)}\left\|f_{i}-\psi\right\|_{2^{*} ; \Omega} \leq C\left\|\nabla\left(f_{i}-\psi\right)\right\|_{2 ; \Omega}
$$

This along with (11.66) shows that $\left\{\left\|f_{i}\right\|_{1,2 ; \Omega}\right\}_{i=1}^{\infty}$ is a bounded sequence. Referring to Remark 11.15 , we know that $W^{1,2}(\Omega)$ is a reflexive banach space and therefore Theorem 8.37 implies that there is a subsequence (denoted by the full sequence) and $f \in W^{1,2}(\Omega)$ such that $f_{i} \rightarrow f$ weakly in $W^{1,2}(\Omega)$. This is equivalent to

$$
\begin{array}{r}
f_{i} \rightarrow f \text { weakly in } L^{2}(\Omega) \\
\nabla f_{i} \rightarrow \nabla f \text { weakly in } L^{2}(\Omega) . \tag{11.68}
\end{array}
$$

Furthermore, it follows from the lower semicontinuity of the norm (note that Theorem 8.35 (iii) is also true with $\|x\|^{2} \leq \lim _{k \rightarrow \infty}\left\|x_{k}\right\|^{2}$ ) and (11.68) that

$$
\int_{\Omega}|\nabla f|^{2} d x \leq \liminf _{i \rightarrow \infty} \int_{\Omega}\left|\nabla f_{i}\right|^{2} d x
$$

To show that $f$ is a valid competitor in (11.65) we need to establish that $f-\psi \in W_{0}^{1,2}(\Omega)$, for then we will have

$$
\int_{\Omega}|\nabla f|^{2} d x=m
$$

thus establishing that the infimum in (11.65) is attained. To show that $f$ assumes the correct boundary values, note that (for a subsequence) $f_{i}-\psi \rightarrow g$ weakly for some $g \in W_{0}^{1,2}(\Omega)$ since $\left\{\left\|f_{i}-\psi\right\|_{1,2 ; \Omega}\right\}$ is a bounded sequence. But $f_{i}-\psi \rightarrow f-\psi$ weakly in $W^{1,2}(\Omega)$ and therefore $f-\psi=g \in W_{0}^{1,2}(\Omega)$.

The next step is to show that $f$ is weakly harmonic. Choose $\varphi \in C_{c}^{\infty}(\Omega)$ and for each real number $t$, let

$$
\begin{aligned}
\alpha(t) & =\int_{\Omega}|\nabla(f+t \varphi)|^{2} d x \\
& =\int_{\Omega}|\nabla f|^{2}+2 t \nabla f \cdot \nabla \varphi+t^{2}|\nabla \varphi|^{2} d x
\end{aligned}
$$

Note that $\alpha$ has a local minimum at $t=0$. Furthermore, referring to Exercise 7, Section 6.2, we see that it is permissible to differentiate under the integral sign to compute $\alpha^{\prime}(t)$. Thus, it follows that

$$
0=\alpha^{\prime}(0)=2 \int_{\Omega} \nabla f \cdot \nabla \varphi d x
$$

which shows that $f$ is weakly harmonic.

## Exercises for Section 11.6

1. Let $f \in C^{2}(\Omega)$. Show that $f$ is harmonic in $\Omega$ if and only if $f$ is weakly harmonic in $\Omega$ (i.e., $\int_{\Omega} f \Delta \varphi d x=0$ for every $\left.\varphi \in C_{c}^{\infty}(\Omega)\right)$. Show also that $f$ is weakly harmonic in $\Omega$ if and only if $\int_{\Omega} \nabla \varphi \cdot \nabla f d x=0$ for every $\varphi \in C_{c}^{\infty}(\Omega)$.
2. Let $f \in W^{1,2}(\Omega)$. Show that $f$ is weakly harmonic (i.e.,

$$
\int_{\Omega} f \Delta \varphi d x=0
$$

for every $\left.\varphi \in C_{c}^{\infty}(\Omega)\right)$ if and only if

$$
\int_{\Omega} \nabla f \cdot \nabla \varphi d x=0
$$

for every test function $\varphi$. Hint: Use Theorem 11.21 along with (11.42).

### 11.7. Regularity of Weakly Harmonic Functions

We will now show that the weak solution found in the previous section is actually a classical $C^{\infty}$ solution.
11.30. THEOREM. If $f \in W_{\mathrm{loc}}^{1,2}(\Omega)$ is weakly harmonic, then $f$ is continuous in $\Omega$ and

$$
f\left(x_{0}\right)=f_{B\left(x_{0}, r\right)} f(y) d y
$$

whenever $\bar{B}\left(x_{0}, r\right) \subset \Omega$.
Proof. Step 1: We will proof first that, for every $x_{0} \in \Omega$, the function

$$
\begin{equation*}
F(r)=\int_{\partial B\left(x_{0}, 1\right)} f(r, z) d \mathcal{H}^{n-1}(z) \tag{11.69}
\end{equation*}
$$

is constant for all $0<r<d\left(x_{0}, \Omega\right)$. Without loss of generality we assume $x_{0}=0$. In order to prove (11.69) we recall that, since $f \in W_{l o c}^{1,2}(\Omega)$ is weakly harmonic, then

$$
\begin{equation*}
\int_{\Omega} f \Delta \varphi d \lambda(x)=0, \quad \text { for all } \varphi \in C_{c}^{\infty}(\Omega) \tag{11.70}
\end{equation*}
$$

We want to choose an appropriate test function $\varphi$ in (11.70). We consider a test function of the form

$$
\begin{equation*}
\varphi(x)=\omega(|x|) \tag{11.71}
\end{equation*}
$$

A direct calculation shows

$$
\begin{aligned}
\frac{\partial \varphi}{\partial x_{i}} & =\omega^{\prime}(|x|) \frac{\partial}{\partial x_{i}}\left(x_{1}^{2}+\ldots+x_{n}^{2}\right)^{\frac{1}{2}} \\
& =\omega^{\prime}(|x|) \frac{1}{2}\left(x_{1}^{2}+\ldots+x_{n}^{2}\right)^{-\frac{1}{2}}\left(2 x_{i}\right) \\
& =\omega^{\prime}(|x|) \frac{x_{i}}{|x|}
\end{aligned}
$$

and hence

$$
\frac{\partial^{2} \varphi}{\partial x_{i}^{2}}=\frac{x_{i}}{|x|}\left(\omega^{\prime \prime}(|x|) \frac{x_{i}}{|x|}\right)+\omega^{\prime}(|x|)\left[\frac{|x|-x_{i} \frac{x_{i}}{|x|}}{|x|^{2}}\right]=\frac{x_{i}^{2}}{|x|^{2}} \omega^{\prime \prime}(|x|)+w^{\prime}(|x|)\left[\frac{|x|^{2}-x_{i}^{2}}{|x|^{3}}\right], i=1,2,3 \ldots
$$

Therefore,

$$
\frac{\partial^{2} \varphi}{\partial x_{i}^{2}}+\ldots+\frac{\partial^{2} \varphi}{\partial x_{n}^{2}}=\omega^{\prime \prime}(|x|)+\frac{n}{|x|} \omega^{\prime}(|x|)-\frac{\omega^{\prime}(|x|)}{|x|}
$$

from which we conclude

$$
\Delta \varphi(x)=\omega^{\prime \prime}(|x|)+\frac{(n-1)}{|x|} \omega^{\prime}(|x|)
$$

Let $r=|x|$. Let $0<t<T<d\left(x_{0}, \partial \Omega\right)$. We choose $w(r)$ such that $w \in C_{c}^{\infty}(t, T)$, and with this $w$ we compute

$$
\begin{aligned}
0 & =\int_{\Omega} f(x) \Delta \varphi(x) d \lambda(x) \\
& =\int_{\Omega} f(x)\left[w^{\prime \prime}(|x|)+\frac{(n-1)}{|x|} w^{\prime}(|x|)\right] d \lambda(x) \\
& =\int_{t}^{T} \int_{\partial B(0,1)} f(r, z)\left[w^{\prime \prime}(r)+\frac{(n-1)}{r} w^{\prime}(r)\right] r^{n-1} d \mathcal{H}^{n-1}(z) d r \\
& =\int_{t}^{T} \int_{\partial B(0,1)} f(r, z)\left[w^{\prime \prime}(r) r^{n-1}+(n-1) r^{n-2} w^{\prime}(r)\right] d \mathcal{H}^{n-1}(z) d r \\
& =\int_{t}^{T} \int_{\partial B(0,1)} f(r, z)\left(r^{n-1} w^{\prime}(r)\right)^{\prime} d \mathcal{H}^{n-1}(z) d r \\
& =\int_{t}^{T} F(r)\left(r^{n-1} w^{\prime}(r)\right)^{\prime} d r
\end{aligned}
$$

We conclude

$$
\begin{equation*}
\int_{t}^{T} F(r)\left(r^{n-1} w^{\prime}(r)\right)^{\prime} d r=0 \tag{11.72}
\end{equation*}
$$

for every test function of the form (11.71), $w \in C_{c}^{\infty}(t, T)$. Given any $\psi \in C_{c}^{\infty}(t, T)$, $\int_{t}^{T} \psi=0$, we now proceed to construct a particular function $w(r)$ such that $\left(r^{n-1} w^{\prime}(r)\right)^{\prime}=\psi$. Indeed, for each real number $r$, define:

$$
y(r)=\int_{t}^{r} \psi(s) d s
$$

and define $\eta$ by

$$
\eta(r)=\int_{t}^{r} \frac{y(s)}{s^{n-1}} d s
$$

Finally, let

$$
w(r)=\eta(r)-\eta(T)
$$

Note that $w \equiv 0$ on $[0, t]$ and $[T, \infty)$. Since $w^{\prime}(r)=\frac{y(r)}{r^{n-1}}$ we obtain from (11.72)

$$
\begin{aligned}
0 & =\int_{t}^{T} F(r)\left(r^{n-1} \frac{y(r)}{r^{n-1}}\right)^{\prime} d r \\
& =\int_{t}^{T} F(r) y^{\prime}(r) d r \\
& =\int_{t}^{T} F(r) \psi(r) d r
\end{aligned}
$$

We have proved that

$$
\begin{equation*}
\int_{t}^{T} F(r) \psi(r) d r=0, \text { for every } \psi \in C_{c}^{\infty}(t, T), \int_{t}^{T} \psi=0 \tag{11.73}
\end{equation*}
$$

We consider the function $F(r)$ as a distribution, say $T_{F}$ (see Definition 10.3). It is clear that (11.73) implies

$$
\int_{t}^{T} F(r) \varphi^{\prime}(r)=0, \text { for every } \varphi \in C_{c}^{\infty}(t, T)
$$

and, from Definition 10.7, this is equivalent to

$$
\begin{equation*}
T_{F}^{\prime}(\varphi)=0, \text { for every } \varphi \in C_{c}^{\infty}(t, T) \tag{11.74}
\end{equation*}
$$

From (11.74) we conclude that $T_{F}^{\prime}$, in the sense of distributions, is 0 and we now appeal to Theorem 10.8 to conclude that the distribution $T_{F}$ is a constant, say $\alpha$. Therefore, we have shown that $F(r)=\alpha\left(x_{0}\right)$ for all $t<r<T$. Since $t$ and $T$ are arbitrary we conclude that $F(r)=\alpha\left(x_{0}\right)$ for all $0<r<d\left(x_{0}, \partial \Omega\right)$, which is (11.69).

Step 2: In this step we show that

$$
\begin{equation*}
f_{B\left(x_{0}, \delta\right)} f(y) d \lambda(y)=C\left(x_{0}, n\right) \tag{11.75}
\end{equation*}
$$

for every $x_{0} \in \Omega$ and every $0<\delta<d\left(x_{0}, \partial \Omega\right)$. Without loss of generality we assume $x_{0}=0$. We fix $0<\delta<d\left(x_{0}, \partial \Omega\right)$. We have, with the help of step 1 ,

$$
\begin{aligned}
\int_{B(0, \delta)} f(x) d \lambda(x) & =\int_{0}^{\delta} \int_{\partial B(0,1)} f(r, z) r^{n-1} d \mathcal{H}^{n-1}(z) d r \\
& =\int_{0}^{\delta} \alpha(0) r^{n-1} d r \\
& =\alpha(0) \frac{\delta^{n}}{n}
\end{aligned}
$$

Hence,

$$
\frac{1}{\delta^{n}} \int_{B\left(x_{0}, \delta\right)} f(x) d \lambda(x)=\frac{\alpha\left(x_{0}\right)}{n}
$$

and therefore,

$$
\begin{equation*}
\frac{1}{\lambda\left(B\left(x_{0}, \delta\right)\right)} \int_{B\left(x_{0}, \delta\right)} f(x) d \lambda(x)=\alpha\left(x_{0}\right) C(n) \tag{11.76}
\end{equation*}
$$

Step 3: Since almost every $x_{0} \in \Omega$ is a Lebesgue point for $f$, we deduce from (11.76) that

$$
f\left(x_{0}\right)=\frac{1}{\lambda\left(B\left(x_{0}, \delta\right)\right)} \int_{B\left(x_{0}, \delta\right)} f(x) d \lambda(x)=\alpha\left(x_{0}\right) C(n)
$$

for $\lambda$-a.e. $x_{0} \in \Omega$ and any ball $\bar{B}\left(x_{0}, \delta\right) \subset \Omega$.
Step 4: $f$ can now be redefined on a set of measure zero in such a way as to ensure its continuity in $\Omega$. Indeed, if $x_{0}$ is not a Lebesgue point and $\bar{B}\left(x_{0}, r\right) \subset \Omega$, we define, with the aid of Step 2,

$$
f\left(x_{0}\right):=\int_{B\left(x_{0}, r\right)} f(y) d \lambda(y)
$$

We left as an exercise to show that, with this definition, $f$ is continuous in $\Omega$.
11.31. Definition. A function $\varphi: B\left(x_{0}, r\right) \rightarrow \mathbb{R}$ is called radial relative to $x_{0}$ if $\varphi$ is constant on $\partial B\left(x_{0}, t\right), t \leq r$.
11.32. Corollary. Suppose $f \in W^{1,2}(\Omega)$ is weakly harmonic and $\bar{B}\left(x_{0}, r\right) \subset$ $\Omega$. If $\varphi \in C_{c}^{\infty}\left(B\left(x_{0}, r\right)\right)$ is nonnegative and radial relative to $x_{0}$ with

$$
\int_{B\left(x_{0}, r\right)} \varphi(x) d x=1
$$

then

$$
\begin{equation*}
f\left(x_{0}\right)=\int_{B\left(x_{0}, r\right)} f(x) \varphi(x) d x \tag{11.77}
\end{equation*}
$$

Proof. For convenience, assume $x_{0}=0$. Since $\varphi$ is radial, note that each superlevel set $\{\varphi>t\}$ is a ball centered at 0 ; let $r(t)$ denote its radius. Then, with $M:=\sup _{B(0, r)} \varphi$ we compute

$$
\begin{aligned}
\int_{B(0, r)} f(x) \varphi(x) d x & =\int_{B(0, r)}\left(f^{+}(x)-f^{-}(x)\right) \varphi(x) d x \\
& =\int_{B(0, r)} \varphi(x) f^{+}(x) d \lambda(x)-\int_{B(0, r)} \varphi(x) f^{-}(x) d \lambda(x) \\
& =\int_{B(0, r)} \varphi(x) d \mu^{+}(x)-\int_{B(0, r)} \varphi(x) d \mu^{-}(x),
\end{aligned}
$$

where the positive measures $\mu^{+}$and $\mu^{-}$are given by

$$
\mu^{+}(E)=\int_{E} f^{+}(x) d \lambda(x) \text { and } \mu^{-}(E)=\int_{E} f^{-}(x) d \lambda(x)
$$

Hence, from Theorem 6.59,

$$
\int_{B(0, r)} f(x) \varphi(x) d x=\int_{0}^{M} \mu^{+}(\{x: \varphi(x)>t\}) d \lambda(t)-\int_{0}^{M} \mu^{-}(\{x: \varphi(x)>t\}) d \lambda(t)
$$

Since $\{\varphi>t\}=B(0, r(t))$, we have

$$
\begin{aligned}
\int_{B(0, r)} f(x) \varphi(x) d x & =\int_{0}^{M} \int_{\{\varphi>t\}} d \mu^{+}(x) d \lambda(t)-\int_{0}^{M} \int_{\{\varphi>t\}} d \mu^{-}(x) d \lambda(t) \\
& =\int_{0}^{M} \int_{\{\varphi>t\}}\left(f^{+}(x)-f^{-}(x)\right) d \lambda(x) d \lambda(t) \\
& =\int_{0}^{M} \int_{B(0, r(t))} f(x) d \lambda(x) d \lambda(t)
\end{aligned}
$$

Appealing now to Theorem 11.30 we conclude

$$
\begin{aligned}
\int_{B(0, r)} f(x) \varphi(x) d x & =\int_{0}^{M} f(0) \lambda[B(0, r(t))] d \lambda(t) \\
& =f(0) \int_{0}^{M} \lambda[\{\varphi>t\}] d t \\
& =f(0) \int_{B(0, r)} \varphi(x) d \lambda(x) \\
& =f(0)
\end{aligned}
$$

We have shown (11.77).
11.33. Theorem. A weakly harmonic function $f \in W^{1,2}(\Omega)$ is of class $C^{\infty}(\Omega)$.

Proof. For each domain $\Omega^{\prime} \subset \subset \Omega$ we will show that $f(x)=f * \varphi_{\varepsilon}(x)$ for $x \in \Omega^{\prime}$. Since $f * \varphi_{\varepsilon} \in C^{\infty}\left(\Omega^{\prime}\right)$, this will suffice to establish our result. As usual, $\varphi_{\varepsilon}(x):=\varepsilon^{-n} \varphi(x / \varepsilon), \varphi \in C_{c}^{\infty}(B(0,1))$ and $\int_{B(0,1)} \varphi(x) d \lambda(x)=1$. We also require
$\varphi$ to be radial with respect to 0 . Finally, we take $\varepsilon$ small enough to ensure that $f * \varphi_{\varepsilon}$ is defined on $\Omega^{\prime}$. We have, for each $x \in \Omega^{\prime}$,
$f * \varphi_{\varepsilon}(x)=\int_{\Omega} \varphi_{\varepsilon}(x-y) f(y) d y=\int_{B(x, \varepsilon)} \varphi_{\varepsilon}(x-y) f(y) d y=\int_{B(0, \varepsilon)} \varphi_{\varepsilon}(y) f(x-y) d y$.
By Exercise 1, Section 11.7, the function $h(y):=f(x-y)$ is weakly harmonic in $B(0, \varepsilon)$. Therefore, since $\varphi_{\varepsilon}$ is radial and $\int_{B(0, \varepsilon)} \varphi_{\varepsilon}(x) d \lambda(x)=1$ we obtain, with the help Corollary 11.32,

$$
f * \varphi_{\varepsilon}(x)=\int_{B(0, \varepsilon)} h(y) \varphi_{\varepsilon}(y) d y=h(0)=f(x)
$$

Theorem 11.29 states that if $\Omega$ is a bounded open set and $\psi \in W^{1,2}(\Omega)$ a given function, then there exists a weakly harmonic function $f \in W^{1,2}(\Omega)$ such that $f-\psi \in W_{0}^{1,2}(\Omega)$. That is, $f$ assumes the same values as $\psi$ on the boundary of $\Omega$ in the sense of Sobolev theory. The previous result shows that $f$ is a classically harmonic function in $\Omega$. If $\psi$ is known to be continuous on the boundary of $\Omega$, a natural question is whether $f$ assumes the values $\psi$ continuously on the boundary. The answer to this is well understood and is the subject of other areas in analysis.

## Exercises for Section 11.7

1. Suppose $f \in W^{1,2}(\Omega)$ is weakly harmonic and let $x \in \Omega^{\prime} \subset \subset \Omega$. Choose $r \leq$ dist $\left(\Omega^{\prime}, \partial \Omega\right)$. Prove that the function $h(y):=f(x-y)$ is weakly harmonic in $B(0, r)$.

## Bibliography

This list includes books and articles that were cited in the text and some additional references that will be useful for further study.
[1] L. Ambrosio, N. Fusco, and D. Pallara. Functions of Bounded Variation and Free Discontinuity Problems. Oxford Mathematical Monographs. The Clarendon Press, Oxford University Press: New York, 2000.
[2] St. Banach and A. Tarski. Sur la décomposition des ensembles de points en parties respectivement congruentes. Fund. Math., 6:244-277, 1924.
[3] T. Bagby, and W. P. Ziemer. Pointwise differentiability and absolute continuity. Trans. Amer. Math. Soc., 194:129-148, 1974.
[4] A. S. Besicovitch. A general form of the covering principle and relative differentiation of additive functions. Proc. Camb. Phil. Soc., 41:103-110, 1945.
[5] A. S. Besicovitch. A general form of the covering principle and relative differentiation of additive functions, II. Proc. Camb. Phil. Soc., 42:1-10, 1946.
[6] B. Bojarski. Remarks on some geometric properties of Sobolev mappings. Functional Analysis and Related Topics, Editor Shozo Koshi, World Scientific, 1991.
[7] J. Bourgain and H. Brezis. On the equation $\operatorname{div} Y=f$ and application to control of phases. J. Amer. Math. Soc., 16(1):393-426, 2002.
[8] L. Carleson. On the connection between Hausdorff measures and capacity. Ark. Mat., 3: 403-406, 1958.
[9] L. Carleson. Selected problems on exceptional sets. Van Nostrand Company, Inc., Priceton, New Jersey, Toronto - London - Melbourne, 1967.
[10] G.-Q. Chen and H. Frid. Extended divergence-measure fields and the Euler equations for gas dynamics. Comm. Math. Phys., 236(2):251-280, 2003.
[11] G.Q. Chen and H. Frid. Divergence-measure fields and hyperbolic conservation laws. Arch. Ration. Mech. Anal., 147(2):89-118, 1999.
[12] G.-Q. Chen and M. Torres. Divergence-measure fields, sets of finite perimeter, and conservation laws. Arch. Ration. Mech. Anal., 175(2):245-267, 2005.
[13] G.-Q. Chen, M. Torres, and W. Ziemer. Gauss-Green theorem for weakly differentiable vector fields, sets of finite perimeter, and balance laws. Communications on Pure and Applied Mathematics, 62, No 2:242-304, 2009.
[14] G.-Q. Chen, M. Torres, and W. P. Ziemer. Measure-theoretical analysis and non-linear conservation laws. Pure and Applied Mathematics Quarterly, 3(3):841-879, 2007.
[15] P. J. Cohen. Set theory and the continuum hypothesis, Benjamin, New York, 1966.
[16] P. J. Cohen. The independence of the continuum hypothesis, Proceedings of the National Academy of Sciences of the United States of America, 50:1143-1148, 1963.
[17] P. J. Cohen. The independence of the continuum hypothesis. II, Proceedings of the National Academy of Sciences of the United States of America, 51:105-110, 1964.
[18] P. J. Cohen. A minimal model for set theory, Bull. Amer. Math. Soc., 69:537540, 1963.
[19] G. Comi and M. Torres. One sided approximations of sets of finite perimeter. Rendiconti Lincei-Matematica e Applicazioni, to appear, 2017.
[20] T. De Pauw. On SBV dual. Indiana Univ. Math. J., 47(1):99-121, 1998.
[21] T. De Pauw. On the exceptional sets of the flux of a bounded vectorfield. J. Math. Pures Appl., 82: 1191-1217, 2003.
[22] T. De Pauw and W. F. Pfeffer. Distributions for which $\operatorname{div} v=f$ has a continuous solution. Comm. Pure Appl. Math., 61(2):230-260, 2008.
[23] T. De Pauw and W. F. Pfeffer. The Gauss-Green theorem and removable sets for PDEs in divergence form. Adv. Math., 183(1):155-182, 2004.
[24] T. De Pauw and M. Torres. On the distributional divergence of vector fields vanishing at infinity. Proceedings of the Royal Society of Edinburgh, 141A:6576, 2011.
[25] L. C. Evans and R. F. Gariepy. Measure Theory and Fine Properties of Functions. Studies in Advanced Mathematics. CRC Press: Boca Raton, FL, 1992.
[26] K. Falconer. Fractal Geometry. Wiley, New York, 1990.
[27] H. Federer. Geometric Measure Theory. Springer-Verlag, New York, Heidelberg, 1969.
[28] H. Federer and W. Ziemer. The Lebesgue set of a function whose distributional derivatives are p-th power summable. Indiana Univ. Math. J., 22:139-158, 1972.
[29] B. Fuglede. Extremal length and functional completion. Acta Mathematica, 98(130):171-219, 1957.
[30] E. Giusti. Minimal Surfaces and Functions of Bounded Variation. With notes by G. H. Williams, Notes on Pure Mathematics, 10, Department of Pure Mathematics, Australian National University, Canberra, 1977.
[31] K. Gödel. The Consistency of the Continuum Hypothesis. Annals of Mathematics Studies, no 3., Princeton University Press, Princeton, N. J., 1940.
[32] W. Gustin. Boxing inequalities. J. Math. Mech., 9:229-239, 1960.
[33] P. Halmos. Naive Set Theory. Van Nostrand, Princeton, N.J., 1960.
[34] G. H. Hardy. Weierstrass's nondifferentiable function. Transactions of the American Mathematical Society, 17(2):301-325, 1916.
[35] F. Maggi. Sets of finite perimeter and geometric variational problems. Cambridge studies in advance mathematics, 2012.
[36] V. G. Maz'ja. Sobolev spaces. Springer Series in Soviet Mathematics. SpringerVerlag, Berlin, 1985. Translated from the Russian by T. O. Shaposhnikova.
[37] N. G. Meyers and W. P. Ziemer. Integral inequalities of Poincare and Wirtinger type for $B V$ functions. Amer. J. Math., 99:1345-1360, 1977.
[38] C. B. Morrey Jr. Functions of several variables and absolute continuity, II. Duke Math. J., 6:187-215, 1940.
[39] A. P. Morse. The behavior of a function on its critical set. Ann. Math., 40:62-70, 1939.
[40] A. P. Morse. Perfect blankets. Trans. Amer. Math. Soc., 6:418-442, 1947.
[41] W. F. Pfeffer. Derivation and integration, volume 140 of Cambridge Tracts in Mathematics, Cambridge University Press, Cambridge, 2001.
[42] N. C. Phuc and M. Torres. Characterizations of the existence and removable singularities of divergence-measure vector fields. Indiana University Mathematics Journal, 57(4):1573-1598, 2008.
[43] N. C. Phuc and M. Torres. Characterizations of signed measures in the dual of $B V$ and related isometric isomorphisms. Annali della Scuola Normale Superiore di Pisa, issue 2 of the XVII volume, series V, 2017.
[44] S. Sacks. Selected problems on exceptional sets. Van Nostrand Company, Inc., Princeton, New Jersey, Toronto - London - Melbourne, 1967.
[45] A. R. Schep. And still one more proof of the Radon-Nikodym theorem. American Mathematical Monthly, 110(6):536-538, 2003.
[46] L. Schwartz. Theorie des Distributions, 2nd edition, Hermann, Paris, 1966.
[47] M. Šilhavý. Divergence-measure fields and Cauchy's stress theorem. Rend. Sem. Mat Padova, 113:15-45, 2005.
[48] E. M. Stein. Singular integrals and differentiability of functions. Princeton University Press, Priceton, New Jersey, 1970.
[49] K. Weierstrass. On continuous functions of a real argument which possess a definite derivative for no value of the argument. Kniglich Preussichen Akademie der Wissenschaften, 110(2):71-74, 1895.
[50] W. P. Ziemer. Weakly differentiable functions, volume 120 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1989.

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[^0]:    ${ }^{1}$ Note the similarity between the technique used in the following argument and the proof of Theorem 4.16, from (4.11) to the end of that proof

[^1]:    In integration theory, functions that differ only on a set of measure zero can be identified as one function. Consequently, with this identification a measurable function determines an equivalence class of functions. This raises the question of whether it is possible to define a measurable function at almost all points in a way that is independent of any representative in the equivalence class. Our investigation of Lebesgue points provides a positive answer to this question.

