# Density convergence in the Breuer-Major theorem for Gaussian stationary sequences 

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Consider a Gaussian stationary sequence with unit variance $X=\left\{X_{k} ; k \in \mathbb{N} \cup\{0\}\right\}$. Assume that the central limit theorem holds for a weighted sum of the form $V_{n}=n^{-1 / 2} \sum_{k=0}^{n-1} f\left(X_{k}\right)$, where $f$ designates a finite sum of Hermite polynomials. Then we prove that the uniform convergence of the density of $V_{n}$ towards the standard Gaussian density also holds true, under a mild additional assumption involving the causal representation of $X$.

Keywords: Breuer-Major theorem; density convergence; Gaussian stationary sequences; Malliavin calculus; moving average representation

## 1. Introduction

Let $X=\left\{X_{k} ; k \in \mathbb{N} \cup\{0\}\right\}$ be a centered Gaussian stationary sequence with unit variance. For all $v \in \mathbb{Z}$, we set $\rho(v)=\mathbf{E}\left[X_{0} X_{|v|}\right]$. Therefore, $\rho(0)=1$ and $|\rho(v)| \leq 1$ for all $v$. Let $\gamma$ be the standard Gaussian probability measure and $f \in L^{2}(\gamma)$ be a fixed deterministic function such that $\mathbf{E}\left[f\left(X_{1}\right)\right]=0$. We expand $f$ in the orthonormal basis of Hermite polynomials $\left\{H_{k} ; k \geq 0\right\}$, which are more specifically defined in Section 2.2. In particular, if $f$ has Hermite rank $d \geq 1$, it admits the following Hermite expansion:

$$
f(x)=\sum_{j=d}^{\infty} a_{j} H_{j}(x)
$$

with $a_{d} \neq 0$. Define $V_{n}=\frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} f\left(X_{k}\right)$. Then the celebrated Breuer-Major theorem (see [3] or Theorem 7.2.4 in [14]) can be written as follows:

Theorem 1.1. Suppose that $\sum_{v \in \mathbb{Z}}|\rho(v)|^{d}<\infty$ and set $\sigma^{2}=\sum_{j=d}^{\infty} j!a_{j}^{2} \sum_{v \in \mathbb{Z}} \rho(v)^{j}$, which is assumed to be in $(0, \infty)$. Then the convergence

$$
\begin{equation*}
V_{n} \xrightarrow{\text { Law }} \mathcal{N}\left(0, \sigma^{2}\right) \tag{1}
\end{equation*}
$$

holds true as $n$ tends to infinity.

We shall be in fact interested in a particular case of Theorem 1.1 for finite linear combinations of Hermite polynomials, which is stated here for convenience.

Corollary 1.2. Consider $2 \leq d \leq q<\infty$ and a family of real numbers $\left\{a_{j} ; j=d, \ldots, q\right\}$. Let $H_{j}$ be the $j$ th order Hermite polynomial, and assume that $\sigma^{2} \in(0, \infty)$, where $\sigma^{2} \equiv$ $\sum_{j=d}^{q} j!a_{j}^{2} \sum_{v \in \mathbb{Z}} \rho(v)^{j}$. Set

$$
\begin{equation*}
V_{n}^{d, q}=\frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} \sum_{j=d}^{q} a_{j} H_{j}\left(X_{k}\right) . \tag{2}
\end{equation*}
$$

Then $V_{n}^{d, q} \xrightarrow{\text { Law }} \mathcal{N}\left(0, \sigma^{2}\right)$ as $n$ tends to infinity. In particular, we have

$$
\begin{equation*}
\lim _{n \rightarrow 0} \mathbf{E}\left[\left(V_{n}^{d, q}\right)^{4}\right]=3 \sigma^{4} \tag{3}
\end{equation*}
$$

Remark 1.3. The relation between Gaussian convergence in law for sequences in a fixed Wiener chaos and behavior of the fourth moment has been extensively studied since the seminal paper [18]. We will need only a small part of the information available on the topic, such as relation (3).

Due to its importance, Breuer-Major theorem has been extended and refined in several directions. Important generalizations can be found in Arcones [1] (multidimensional case), Chambers and Slud [6] and Giraitis and Surgailis [7]. A proof of Theorem 1.1 using a combination of Stein's method with Malliavin calculus was given by Nourdin, Peccati and Podolskij in [15], where one can find explicit bounds in the total variation and Wasserstein distances. We refer the reader to the monograph by Nourdin and Peccati [13] for a more detailed account on this topic.

We shall mainly be concerned here by convergences of densities, and here again the relationship between fourth moment behavior and various type of convergences of random variables in a fixed Wiener chaos have been thoroughly studied in the recent past. The interested reader is referred to [14] for further details, but we will use here the following recent criterion (see [12], Corollary 1.2, and [11], Corollary 4.6).

Theorem 1.4. Let $\left\{F_{n} ; n \in \mathbb{N}\right\}$ be a sequence of random variables belonging to a fixed chaos $\mathcal{H}_{q}$ with $q \geq 2$. Suppose $\mathbf{E}\left[F_{n}^{2}\right]=1$ and $\lim _{n \rightarrow \infty} \mathbf{E}\left[F_{n}^{4}\right]=3$. Let $p_{F_{n}}$ be the density of the random variable $F_{n}$ and let $\phi(x)=(2 \pi)^{-1 / 2} \exp \left(-|x|^{2} / 2\right)$ be the density of the standard Gaussian distribution on $\mathbb{R}$.
(i) Suppose that for some $\varepsilon>0$,

$$
\sup _{n} \mathbf{E}\left[\left\|D F_{n}\right\|^{-4-\varepsilon}\right]<\infty
$$

Then, there exists a constant $c$ such that for all $n \geq 1$,

$$
\sup _{x \in \mathbb{R}}\left|p_{F_{n}}(x)-\phi(x)\right| \leq c \sqrt{\mathbf{E}\left[F_{n}^{4}\right]-3} .
$$

(ii) Suppose that for all $p \geq 1$,

$$
\sup _{n} \mathbf{E}\left[\left\|D F_{n}\right\|^{-p}\right]<\infty
$$

Then, for any $m \geq 0$, there exists a constant $c_{m}$ such that for all $n \geq 1$,

$$
\sup _{x \in \mathbb{R}}\left|p_{F_{n}}^{(m)}(x)-\phi^{(m)}(x)\right| \leq c_{m} \sqrt{\mathbf{E}\left[F_{n}^{4}\right]-3}
$$

The goal of the current paper is to apply the criterion given by Theorem 1.4 in order to get convergence of density in the landmark of Breuer-Major theorem. In order to do this, we need a uniform estimate on the negative moments of the Malliavin derivative of the sequence, and this is the contents of our main result.

Theorem 1.5. Let $X$ be a Gaussian stationary sequence whose spectral density $f_{\rho}$ satisfies $\log \left(f_{\rho}\right) \in L^{1}([-\pi, \pi])$ (see Hypothesis 2.1 and the examples in the next section). Let $V_{n}^{d, q}$ be the random variable defined by (2), and assume the hypothesis of Corollary 1.2 to be satisfied. Then for any $p \geq 1$, there exists $n_{0}$ such that

$$
\begin{equation*}
\sup _{n \geq n_{0}} \mathbf{E}\left[\left\|D V_{n}^{d, q}\right\|^{-p}\right]<\infty \tag{4}
\end{equation*}
$$

In the case of a fixed Wiener chaos, we can obtain the following consequence.
Corollary 1.6. Under the conditions of Theorem 1.5 , if $q=d$, and we define $F_{n}=V_{n}^{d, d} / \sigma_{n}$, where $\sigma_{n}^{2}=\mathbf{E}\left[\left(V_{n}^{d, d}\right)^{2}\right]$, then, for all $m \geq 0$ there exists an $n_{0}$ (depending on $m$ ) such that

$$
\sup _{n \geq n_{0}} \sup _{x \in \mathbb{R}}\left|p_{F_{n}}^{(m)}(x)-\phi^{(m)}(x)\right| \leq c_{m} \sqrt{\mathbf{E}\left[F_{n}^{4}\right]-3} .
$$

In the case $q \neq d$, Theorem 1.4 cannot be applied. In the reference [11] one can find results on the uniform convergence of density for general random variables similar to those stated in Theorem 1.4, but they require a uniform lower bound for the negative moments of the random variable $\left|\left\langle D F_{n},-D L^{-1} F_{n}\right\rangle_{\mathfrak{H}}\right|$, and our approach does not seem to work in this case because it is not clear how to express $\left\langle D F_{n},-D L^{-1} F_{n}\right\rangle_{\mathfrak{H}}$ as a sum of squares. Nevertheless, condition (4) allows us to derive the uniform convergence of the densities and their derivatives from a general result proved below (see Proposition 2.6) although in this case we have no information about the rate of convergence.

Corollary 1.7. Under the conditions of Theorem 1.5 , if we define $F_{n}=V_{n}^{d, q} / \sigma_{n}$, where $\sigma_{n}^{2}=$ $\mathbf{E}\left[\left(V_{n}^{d, q}\right)^{2}\right]$, then, for all $m \geq 0$ we have

$$
\lim _{n \rightarrow \infty} \sup _{x \in \mathbb{R}}\left|p_{F_{n}}^{(m)}(x)-\phi^{(m)}(x)\right|=0
$$

Notice that a particular case of Theorem 1.5 has been established in [12], for $q=2$ and $X_{k}=$ $B_{k+1}-B_{k}$ for a fractional Brownian motion $B$ with Hurst parameter $H \in(0,1)$. The proof of the existence of negative moments for $\left\|D F_{n}\right\|$ there is based on the Volterra representation of $B$, which leads to long computations. In comparison our current Theorem 1.5 is more general, since it is valid for a wide class of Gaussian stationary sequences. Its proof is also significantly simplified. These are achieved by the introduction of two new ingredients in the proof, namely:

- A general formula to compute conditional expectations for random variables of the form $H_{q}\left(X_{k}\right)$.
- Related to the previous item, we heavily resort to the causal representation of $X_{k}$, which is particularly convenient in order to compute conditional expectations.

Here is how our paper is structured: we give some preliminary results concerning Gaussian stationary sequences and related Malliavin calculus in Section 2. We then prove our main Theorem 1.5 in Section 3.

## 2. Preliminaries

This section is devoted to some preliminaries on causal or moving average representations for Gaussian stationary sequences, as well as Malliavin calculus tools which will be used in the sequel.

### 2.1. Moving average representation

The classical results on time series presented here are borrowed from [2,4], to which we refer for further details. Start from our Gaussian stationary sequence $\left\{X_{k} ; k \in \mathbb{N} \cup\{0\}\right\}$ with covariance function $\rho$. We will work under the following assumptions:

Hypothesis 2.1. We suppose that $\rho$ admits a spectral density $f_{\rho}$ such that $\log \left(f_{\rho}\right) \in$ $L^{1}([-\pi, \pi])$.

Condition $\log \left(f_{\rho}\right) \in L^{1}([-\pi, \pi])$ is referred to as purely nondeterministic property in the literature. The interest of dealing with purely nondeterministic sequences is that they admit a so-called causal representation which is particularly convenient for conditional expectation computations. Let us state a result in this direction, which is taken from [4], Chapter 5.

Proposition 2.2. Let $X$ be a Gaussian stationary sequence satisfying Hypothesis 2.1. Then for each $k \in \mathbb{N} \cup\{0\}$ the random variable $X_{k}$ can be decomposed as

$$
\begin{equation*}
X_{k}=\sum_{j \geq 0} \psi_{j} w_{k-j}, \tag{5}
\end{equation*}
$$

where $\left(w_{k}\right)_{k \in \mathbb{Z}}$ is a discrete Gaussian white noise and the coefficients $\psi_{j}$ are deterministic. We can always choose the white noise and the coefficients in such a way that $\psi_{0}>0$.

Proof. The existence of the causal moving average representation is a classical result that can be found, for instance, in [4], Theorem 7.5.2.

Notice that from the expansion (5) we easily obtain:

$$
\begin{equation*}
\rho\left(k_{1}-k_{2}\right)=\sum_{l=-\infty}^{k_{1} \wedge k_{2}} \psi_{k_{1}-l} \psi_{k_{2}-l}, \quad \text { for all } k_{1}, k_{2} \in \mathbb{N} \tag{6}
\end{equation*}
$$

and this relation will be used in the proof of our main theorem.
Let us now turn to examples for which our standing assumptions of Hypothesis 2.1 are met. The following proposition provides two typical and classical cases for which a spectral density exists and satisfies some integrability properties.

Proposition 2.3. Let $\rho$ be the covariance function of $X$. We have the following statements.
(i) If $\rho \in \ell^{1}$, then the spectral density $f_{\rho}$ exists and is a nonnegative bounded function defined on $[-\pi, \pi]$.
(ii) Suppose that $\left\{\rho(k) k^{\alpha}, k>0\right\}$ is positive and it is normalized slowly varying at infinity for some $\alpha \in(0,1)$. That is, for every $\delta>0$, for sufficiently large $k, \rho(k) k^{\alpha+\delta}$ is increasing and $\rho(k) k^{\alpha-\delta}$ is decreasing. Then the spectral density exists and satisfies $\lim _{\lambda \rightarrow 0}|\lambda|^{1-\alpha} f_{\rho}(\lambda)=c_{f}$ for some constant $c_{f}>0$ (see [9]).

We now give two specific and important examples which satisfy Hypothesis 2.1.
Example 2.4. The so-called autoregressive fractionally integrated moving-average (ARFIMA) processes are introduced in [8] and [10]. Denote by $B$ the one lag backward operator ( $B X_{k}=$ $\left.X_{k-1}\right)$. Let $\phi(z)$ and $\theta(z)$ be two polynomials which have no common zeros and such that the zeros of $\phi$ lie outside the closed unit disk $\{z,|z| \leq 1\}$. Suppose that $X_{k}$ is given by

$$
\begin{equation*}
\phi(B) X_{k}=(\operatorname{Id}-B)^{-d} \theta(B) w_{k}, \tag{7}
\end{equation*}
$$

where $-1<d<1 / 2$, and where the operator $(\operatorname{Id}-B)^{-d}$ is defined by:

$$
(\operatorname{Id}-B)^{-d}=\sum_{j=1}^{\infty} \eta_{j} B^{j} \quad \text { with } \eta_{j}=\frac{\Gamma(d+j)}{\Gamma(j+1) \Gamma(d)}
$$

Also notice that in (7) the sequence $\left(w_{k}\right)_{k \in \mathbb{Z}}$ is a discrete Gaussian white noise. It is well-known (see [19], Theorem 3.4 and equation (3.19)) that under the above conditions, $\left\{X_{k}, k \in \mathbb{N}\right\}$ admits a spectral density whose exact expression is:

$$
f(\lambda)=\frac{1}{2 \pi}\left[2 \sin \frac{\lambda}{2}\right]^{-2 d} \frac{\left|\theta\left(\mathrm{e}^{-\mathrm{i} \lambda}\right)\right|^{2}}{\left|\phi\left(\mathrm{e}^{-\mathrm{i} \lambda}\right)\right|^{2}}
$$

It is thus readily checked that Hypothesis 2.1 is satisfied, and hence $X_{k}$ has a causal representation.

Example 2.5. Our second example is the fractional Gaussian noise. Let $\left\{B_{t}, t \geq 0\right\}$ be a fractional Brownian motion of Hurst parameter $H \in(0,1)$. Then $\left\{X_{k}=B_{k+1}-B_{k}, k \in \mathbb{N} \cup\{0\}\right\}$ is a stationary Gaussian process with correlation

$$
\rho(k)=\frac{1}{2}\left[(k+1)^{2 H}-2 k^{2 H}+(k-1)^{2 H}\right] .
$$

Its spectral density (see, e.g., [2], equation (2.17)) is:

$$
f(\lambda)=\frac{1}{2 \pi} \sum_{k=-\infty}^{\infty} \rho(|k|) \mathrm{e}^{\mathrm{i} k \lambda}=2 c_{f}(1-\cos (\lambda)) \sum_{j=-\infty}^{\infty}|2 \pi j+\lambda|^{-2 H-1}, \quad \lambda \in[-\pi, \pi],
$$

where $c_{f}=(2 \pi)^{-1} \sin (\pi H) \Gamma(2 H+1)$. If $H \leq 1 / 2$, it is clear that $\sum_{k=-\infty}^{\infty}|\rho(|k|)|<\infty$. This implies

$$
\sup _{\lambda \in[-\pi, \pi]}|f(\lambda)|<\infty .
$$

If $1 / 2<H<1$, then

$$
0 \leq f(\lambda) \leq 2 c_{f}(1-\cos (\lambda))|\lambda|^{-2 H-1}+2 c_{f} \sum_{j \neq 0}|2 \pi j+\lambda|^{-2 H-1}, \quad \lambda \in[-\pi, \pi] .
$$

The first term is in $L^{1}$ since $H<1$. When $j \neq 0, \int_{-\pi}^{\pi}|2 \pi j+\lambda|^{-2 H-1} \mathrm{~d} \lambda \leq C j^{-2 H}$ for some positive constant $C$. Thus $\int_{-\pi}^{\pi} \sum_{j \neq 0}|2 \pi j+\lambda|^{-2 H-1} \mathrm{~d} \lambda<\infty$, owing to the fact that $H>1 / 2$. Therefore, we have $f \in L^{1}$. Summarizing we have $f \in L^{1}$ for all $H \in(0,1)$. This also implies $\log ^{+} f(\lambda) \in L^{1}$. To see $\log ^{-} f(\lambda) \in L^{1}$, we notice that

$$
f(\lambda) \geq 2 c_{f}(1-\cos (\lambda))|\lambda|^{-2 H-1}
$$

So $\log ^{-} f(\lambda) \leq C+\left|\log \left[(1-\cos (\lambda))|\lambda|^{-2 H-1}\right]\right|$ which is in $L^{1}$. In conclusion, the sequence $X$ satisfies Hypothesis 2.1.

### 2.2. Malliavin calculus

We start by briefly recalling some basic notation and results connected to Gaussian analysis and Malliavin calculus. The reader is referred to $[14,17]$ for details or missing proofs.

### 2.2.1. Wiener space and generalities

Let $\mathfrak{H}$ be a real separable Hilbert space with inner product $\langle\cdot, \cdot\rangle_{\mathfrak{H}}$. The norm of $\mathfrak{H}$ will be denoted by $\|\cdot\|=\|\cdot\|_{\mathfrak{H}}$. Recall that we call isonormal Gaussian process over $\mathfrak{H}$ any centered Gaussian family $W=\{W(h): h \in \mathfrak{H}\}$, defined on a probability space $(\Omega, \mathcal{F}, \mathbf{P})$ and such that $\mathbf{E}[W(h) W(g)]=\langle h, g\rangle_{\mathfrak{H}}$ for every $h, g \in \mathfrak{H}$.

In our application the underlying Gaussian family will be a discrete Gaussian white noise $\left(w_{k}\right)_{k \in \mathbb{Z}}$. The space $\mathfrak{H}$ is given here by $\mathfrak{H}=\ell^{2}(\mathbb{Z})$ (the space of square integrable sequences
indexed by $\mathbb{Z}$ ) equipped with its natural inner product. Set $\left\{\varepsilon^{j} ; j \in \mathbb{Z}\right\}$ for the canonical basis of $\ell^{2}(\mathbb{Z})$, that is $\varepsilon_{k}^{j}=\delta_{j}(k)$. We thus identify $w_{j}$ with $W\left(\varepsilon^{j}\right)$. Assume from now on that our underlying $\sigma$-algebra $\mathcal{F}$ is generated by $W$.

For any integer $q \in \mathbb{N} \cup\{0\}$, we denote by $\mathcal{H}_{q}$ the $q$ th Wiener chaos of $W$. We recall that $\mathcal{H}_{0}$ is simply $\mathbb{R}$ whereas, for any $q \geq 1, \mathcal{H}_{q}$ is the closed linear subspace of $L^{2}(\Omega)$ generated by the family of random variables $\left\{H_{q}(W(h)), h \in \mathfrak{H},\|h\|_{\mathfrak{H}}=1\right\}$, with $H_{q}$ the $q$ th Hermite polynomial given by

$$
\begin{equation*}
H_{q}(x)=(-1)^{q} \mathrm{e}^{x^{2} / 2} \frac{\mathrm{~d}^{q}}{\mathrm{~d} x^{q}}\left(\mathrm{e}^{-x^{2} / 2}\right) \tag{8}
\end{equation*}
$$

Let $\mathcal{S}$ be the set of all cylindrical random variables of the form

$$
F=g\left(W\left(h_{1}\right), \ldots, W\left(h_{n}\right)\right),
$$

where $n \geq 1, h_{i} \in \mathfrak{H}$, and $g$ is infinitely differentiable such that all its partial derivatives have polynomial growth. The Malliavin derivative of $F$ is the element of $L^{2}(\Omega ; \mathfrak{H})$ defined by

$$
D F=\sum_{i=1}^{n} \frac{\partial g}{\partial x_{i}}\left(W\left(h_{1}\right), \ldots, W\left(h_{n}\right)\right) h_{i}
$$

By iteration, for every $m \geq 2$, we define the $m$ th derivative $D^{m} F$. This is an element of $L^{2}\left(\Omega ; \mathfrak{H}^{\odot m}\right)$, where $\mathfrak{H}^{\odot m}$ designates the symmetric $m$ th tensor product of $\mathfrak{H}$. For $m \geq 1$ and $p \geq 1, \mathbb{D}^{m, p}$ denote the closure of $\mathcal{S}$ with respect to the norm $\|\cdot\|_{m, p}$ defined by

$$
\|F\|_{m, p}^{p}=\mathbf{E}\left[|F|^{p}\right]+\sum_{j=1}^{m} \mathbf{E}\left[\left\|D^{j} F\right\|_{\mathfrak{H}^{\otimes j}}^{p}\right]
$$

Set $\mathbb{D}^{\infty}=\bigcap_{m, p} \mathbb{D}^{m, p}$. One can then extend the definition of $D^{m}$ to $\mathbb{D}^{m, p}$. When $m=1$, one simply write $D$ instead of $D^{1}$. As a consequence of the hypercontractivity property of the OrnsteinUhlenbeck semigroup (see, e.g., [14], Theorem 2.7.2), all the $\|\cdot\|_{m, p}$-norms are equivalent in any finite sum of Wiener chaoses.

Finally, let us recall that the Malliavin derivative $D$ satisfies the following chain rule: if $\varphi: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is in $\mathcal{C}_{b}^{1}$ (i.e., belongs to the set of continuously differentiable functions with a bounded derivative) and if $\left\{F_{i}\right\}_{i=1, \ldots, n}$ is a vector of elements of $\mathbb{D}^{1,2}$, then $\varphi\left(F_{1}, \ldots, F_{n}\right) \in \mathbb{D}^{1,2}$ and

$$
\begin{equation*}
D \varphi\left(F_{1}, \ldots, F_{n}\right)=\sum_{i=1}^{n} \frac{\partial \varphi}{\partial x_{i}}\left(F_{1}, \ldots, F_{n}\right) D F_{i} \tag{9}
\end{equation*}
$$

### 2.3. Convergence of densities

Suppose that $F$ is a random variable in $\mathbb{D}^{\infty}$ such that $\mathbf{E}\left[\|D F\|^{-p}\right]<\infty$ for all $p \geq 1$. Then, we know that $F$ has an infinitely differentiable density and there are explicit formulas for the density
and its derivatives (see [17], Proposition 2.1.5). Using this result, we can establish the following criteria for convergence of densities for random variables in a finite sum of Wiener chaoses.

Proposition 2.6. Let $\left\{F_{n} ; n \in \mathbb{N}\right\}$ be a sequence of random variables belonging to a finite sum of Wiener chaoses $\bigoplus_{k=1}^{M} \mathcal{H}_{k}$, which converges in law to a nonzero random variable $F_{\infty}$. Suppose that, for all $p \geq 1$

$$
\begin{equation*}
\sup _{n} \mathbf{E}\left[\left\|D F_{n}\right\|^{-p}\right]<\infty \tag{10}
\end{equation*}
$$

Then, for all $m \geq 0$ the derivative $p_{n}^{(m)}$ of the density of $F_{n}$, converges uniformly and in $L^{p}(\mathbb{R})$ for all $p \geq 1$ to the corresponding derivative of the density of $F_{\infty}$.

Proof. First, notice that by condition (10), the random variable $F_{n}$ has an infinitely differentiable density $p_{n}$, whose derivatives can be expressed by

$$
\begin{equation*}
p_{n}^{(m)}(x)=\mathbf{E}\left[\mathbf{1}_{\{F>x\}} G_{n}^{(m)}\right], \tag{11}
\end{equation*}
$$

where the random variables $G_{n}^{(m)}$ are defined recursively by $G_{n}^{(0)}=\delta\left(\frac{D F_{n}}{\left\|D F_{n}\right\|_{H}^{2}}\right)$ and

$$
G_{n}^{(m)}=-\delta\left(\frac{G_{n}^{(m-1)} D F_{n}}{\left\|D F_{n}\right\|_{H}^{2}}\right),
$$

for any $m \geq 1$. From this formula, it follows that the derivatives $p_{n}^{(m)}$ are uniformly bounded and also uniformly bounded in $L^{p}(\mathbb{R})$ for all $p \geq 1$. Indeed, by [16], Lemma 2.4, we have $\sup _{n} \mathbf{E}\left[\left|F_{n}\right|^{r}\right]<\infty$ for all $r \geq 1$. This uniform bound on the moments, together with the equivalence of the $\|\cdot\|_{m, p}$ norms in any finite sum of Wiener chaoses and condition (10) imply that $\sup _{n}\left\|G_{n}^{(m)}\right\|_{L^{p}(\Omega)}=c_{m, p}<\infty$ for all $m \geq 0$. Then, we can write from (11)

$$
\begin{equation*}
\sup _{n} \sup _{x}\left|p_{n}^{(m)}(x)\right| \leq \sup _{n} \mathbf{E}\left[\left|G_{n}^{(m)}\right|\right]=c_{m, 1}<\infty \tag{12}
\end{equation*}
$$

and, using the fact that $\mathbf{E}\left[G_{n}^{(m)}\right]=0$, we get:

$$
\begin{align*}
\sup _{n}\left|p_{n}^{(m)}(x)\right| & \leq \sup _{n}\left(\mathbf{P}\left(\left|F_{n}\right|>|x|\right) \mathbf{E}\left[\left|G_{n}^{(m)}\right|^{2}\right]\right)^{1 / 2} \\
& \leq c_{m, 2} \sup _{n} \sqrt{\mathbf{P}\left(\left|F_{n}\right|>|x|\right)}  \tag{13}\\
& \leq c_{m, 2} \sup _{n} \mathbf{E}\left[\left|F_{n}\right|^{q}\right]^{1 / 2}|x|^{-q / 2} \leq c|x|^{-q / 2},
\end{align*}
$$

for all $q \geq 1$ and for some constant $c$ depending on $q$ and $m$.
By [16], Theorem 3.1, the sequence $F_{n}$ converges in total variation to $F_{\infty}$, that is, the densities $p_{n}$ converge in $L^{1}(\mathbb{R})$ to the density $p_{\infty}$ of $F_{\infty}$. The boundedness in $L^{p}(\mathbb{R})$ and the uniform bound of $p_{n}$ imply that this convergence holds in $L^{p}(\mathbb{R})$ for any $p \geq 1$.

On the other hand, the estimates (12) and (13) imply that for any $m \geq 1$ and any $p \geq 1$, the sequence $p_{n}^{(m)}$ is uniformly bounded in $L^{p}(\mathbb{R})$. Therefore, for any $m \geq 1$ and any $p \geq 1$ the sequence $p_{n}^{(m)}$ is relatively compact in $L^{p}(\mathbb{R})$. Suppose that a subsequence $\left\{p_{n_{k}}^{(m)}, k \geq 1\right\}$ converges in $L^{p}(\mathbb{R})$ to some limit $\tilde{p}_{\infty}^{(m)}$. This limit must coincide with the $m$ th derivative (in the distribution sense) of $p_{\infty}$, and, therefore, it is unique. This implies that for any $m \geq 1$ and any $p \geq 1$ the sequence $p_{n}^{(m)}$ converges in $L^{p}(\mathbb{R})$ to the $m$ th derivative of $p_{\infty}$.

Finally, the uniform convergence is also easy to establish from the convergence of the densities in $L^{p}(\mathbb{R})$ for all $p \geq 1$.

### 2.4. A key lemma

Our future computations will heavily rely on an efficient way to compute conditional expectations. Towards this aim, we state here some general results. Let us start with a decomposition for Hermite polynomials.

Lemma 2.7. For any $q \geq 1$, let $H_{q}$ be the polynomial defined by relation (8). Consider $y, z \in \mathbb{R}$ and two real parameters $a, b$ with $a^{2}+b^{2}=1$. Then the following relation holds true:

$$
\begin{equation*}
H_{q}(a y+b z)=\sum_{\ell=0}^{q}\binom{q}{\ell} a^{q-\ell} b^{\ell} H_{q-\ell}(y) H_{\ell}(z) \tag{14}
\end{equation*}
$$

Proof. By the definition of the Hermite polynomials, we have

$$
\begin{equation*}
\mathrm{e}^{a t y-(a t)^{2} / 2}=\sum_{i=0}^{\infty} \frac{(a t)^{i}}{i!} H_{i}(y) \quad \text { and } \quad \mathrm{e}^{t b z-(b t)^{2} / 2}=\sum_{j=0}^{\infty} \frac{(b t)^{j}}{j!} H_{j}(z) \tag{15}
\end{equation*}
$$

In the same way, we also obtain

$$
\begin{equation*}
\mathrm{e}^{t(a y+b z)-t^{2} / 2}=\sum_{q=0}^{\infty} \frac{t^{q}}{q!} H_{q}(a y+b z) \tag{16}
\end{equation*}
$$

Since $a^{2}+b^{2}=1$, we obviously have $\mathrm{e}^{a t y-(a t)^{2} / 2} \mathrm{e}^{t b z-(b t)^{2} / 2}=\mathrm{e}^{t(a y+b z)-t^{2} / 2}$. Thus, multiplying the right-hand sides of both identities in (15) we recover the right-hand side of (16), namely:

$$
\sum_{q=0}^{\infty} \frac{t^{q}}{q!} H_{q}(a y+b z)=\sum_{i=0}^{\infty} \frac{(a t)^{i}}{i!} H_{i}(y) \sum_{j=0}^{\infty} \frac{(b t)^{j}}{j!} H_{j}(z)
$$

which easily yields the desired identity (14).
With this preliminary result in hand, we are ready to state our result on conditional expectations.

Proposition 2.8. Let $Y$ and $Z$ be two centered Gaussian random variables such that $Y$ is measurable with respect to a $\sigma$-algebra $\mathcal{G} \subset \mathcal{F}$ and $Z$ is independent of $\mathcal{G}$. Assume that $\mathbf{E}\left[Y^{2}\right]=\mathbf{E}\left[Z^{2}\right]=1$. Then for any $q \geq 1$, and real parameters $a, b$ such that $a^{2}+b^{2}=1$, we have:

$$
\begin{equation*}
\mathbf{E}\left[H_{q}(a Y+b Z) \mid \mathcal{G}\right]=a^{q} H_{q}(Y) \tag{17}
\end{equation*}
$$

Proof. Apply identity (14) in order to decompose $H_{q}(a Y+b Z)$. Then identity (17) follows easily from the fact that $Y$ is $\mathcal{G}$-measurable, $Z$ is independent from $\mathcal{G}$ and Hermite polynomials have 0 mean under a centered Gaussian measure except for $H_{0} \equiv 1$.

### 2.5. Carbery-Wright inequality

In the proof of Theorem 1.5, we will make use of the following inequality due to Carbery and Wright [5], Theorem 8, which is recalled here for convenience.

Proposition 2.9. Let $X=\left(X_{1}, \ldots, X_{n}\right)$ be a Gaussian random vector in $\mathbb{R}^{n}$ and $Q: \mathbb{R}^{n} \rightarrow \mathbb{R}$ a polynomial of degree at most $m$. Then there is a universal constant $c>0$ such that:

$$
\begin{equation*}
\mathbf{E}\left[\left|Q\left(X_{1}, \ldots, X_{n}\right)\right|\right]^{1 / m} \mathbf{P}\left(\left|Q\left(X_{1}, \ldots, X_{n}\right)\right| \leq x\right) \leq c m x^{1 / m}, \quad \text { for all } x>0 \tag{18}
\end{equation*}
$$

## 3. Proof of Theorem 1.5

In this section, we will prove our main result, which amounts to show the inequality (4). This will be done into several steps.
Step 1: Computation of the Malliavin norm. Recall that $V_{n}^{d, q}$ is defined by relation (2), and set for the moment $f=\sum_{j=d}^{q} a_{j} H_{j}$. Invoking relation (9), plus the fact that $D w_{j}=\varepsilon^{j}$ with the notation of Section 2.2.1, we get:

$$
\begin{equation*}
D V_{n}^{d, q}=\frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} f^{\prime}\left(X_{k}\right)\left(\sum_{j \geq 0} \psi_{j} \varepsilon^{k-j}\right)=\frac{1}{\sqrt{n}} \sum_{l=-\infty}^{n-1}\left(\sum_{k=l^{+}}^{n-1} \psi_{k-l} f^{\prime}\left(X_{k}\right)\right) \varepsilon^{l} \tag{19}
\end{equation*}
$$

where $l^{+}=\max \{l, 0\}$. Invoking relation (6), it is thus readily checked that:

$$
\left\langle D V_{n}^{d, q}, D V_{n}^{d, q}\right\rangle_{\mathfrak{H}}=\frac{1}{n} \sum_{k_{1}, k_{2}=0}^{n-1} f^{\prime}\left(X_{k_{1}}\right) \rho\left(k_{1}-k_{2}\right) f^{\prime}\left(X_{k_{2}}\right)
$$

where we recall that $\rho$ is the covariance function of the Gaussian stationary sequence $\left\{X_{k} ; k \geq 0\right\}$. This is consistent with the expression found in [14], Chapter 5. However, in order to write the above expression as sum of some squares, we will start directly from expression (19). Since
$\left\{\varepsilon^{l} ; l \in \mathbb{Z}\right\}$ is an orthonormal basis of $\ell^{2}(\mathbb{Z})$ we obtain:

$$
\left\langle D V_{n}^{d, q}, D V_{n}^{d, q}\right\rangle_{\mathfrak{H}}=\frac{1}{n} \sum_{\ell=-\infty}^{n-1}\left(\sum_{k=\ell^{+}}^{n-1} \psi_{k-\ell} f^{\prime}\left(X_{k}\right)\right)^{2}
$$

Rearranging terms (namely, change $k-\ell$ to $k$ and then $n-\ell-1$ to $m$ ), we end up with:

$$
\begin{aligned}
\left\langle D V_{n}^{d, q}, D V_{n}^{d, q}\right\rangle_{\mathfrak{H}} & \geq \frac{1}{n} \sum_{\ell=0}^{n-1}\left(\sum_{k=0}^{n-\ell-1} f^{\prime}\left(X_{\ell+k}\right) \psi_{k}\right)^{2} \\
& =\frac{1}{n} \sum_{m=0}^{n-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{n-1-(m-k)}\right) \psi_{k}\right)^{2} \equiv A_{n}
\end{aligned}
$$

As a last preliminary step we resort to the fact that $X=\left\{X_{k} ; k \in \mathbb{N} \cup\{0\}\right\}$ is a Gaussian stationary sequence, which allows to assert that $A_{n}$ is identical in law to $B_{n}$ with

$$
B_{n}:=\frac{1}{n} \sum_{m=0}^{n-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{m-k}\right) \psi_{k}\right)^{2}=\frac{1}{n} \sum_{m=0}^{n-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k}\right)^{2}
$$

We will now bound the negative moments of $B_{n}$.
Step 2: Block decomposition. We now wish to apply the Carbery-Wright inequality (18) in order to get bounds for negative moments of $B_{n}$. However, relation (18) only applies to moments of small order, and this is why we proceed to a decomposition of $B_{n}$ into smaller blocks.

Fix thus an integer $N \geq 1$ and let $M=[n / N]$ be the integer part of $n / N$. Then $n \geq N M$ and as a consequence,

$$
B_{n}=\frac{1}{n} \sum_{m=0}^{n-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k}\right)^{2} \geq \frac{1}{n} \sum_{i=0}^{N-1} \sum_{m=i M}^{(i+1) M-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k}\right)^{2}
$$

For $i=0, \ldots, N-1$, define

$$
B_{n}^{i}=\frac{1}{n} \sum_{m=i M}^{(i+1) M-1}\left(\sum_{k=0}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k}\right)^{2}
$$

so that $B_{n} \geq \sum_{i=0}^{N-1} B_{n}^{i}$. Then it is readily checked that:

$$
\begin{equation*}
\left(B_{n}\right)^{-p / 2} \leq \prod_{i=0}^{N-1}\left(B_{n}^{i}\right)^{-p /(2 N)} \tag{20}
\end{equation*}
$$

Recall once again the representation of the sequence $X$ in (5), and denote by $\mathcal{F}_{k}$ the filtration generated by $\left\{w_{\ell}: \ell<k\right\}$. Then starting from (20) we obtain:

$$
\begin{align*}
\mathbf{E}\left[\left(B_{n}\right)^{-p / 2}\right] & \leq \mathbf{E}\left[\prod_{i=0}^{N-1}\left(B_{n}^{i}\right)^{-p /(2 N)}\right]  \tag{21}\\
& =\mathbf{E}\left[\mathbf{E}\left[\left(B_{n}^{N-1}\right)^{-p /(2 N)} \mid \mathcal{F}_{(N-1) M}\right] \prod_{i=0}^{N-2}\left(B_{n}^{i}\right)^{-p /(2 N)}\right] .
\end{align*}
$$

Step 3: Application of Carbery-Wright. Let us go back to the particular situation of $f=$ $\sum_{j=d}^{q} a_{j} H_{j}$, which means in particular that $f^{\prime}=\sum_{j=d}^{q} j a_{j} H_{j-1}$. We are now in a position to apply a conditional version of inequality (18) to the block $\left(B_{n}^{N-1}\right)^{-p /(2 N)}$ in (21). First, we notice

$$
\begin{equation*}
\mathbf{E}\left[\left(B_{n}^{N-1}\right)^{-p /(2 N)} \mid \mathcal{F}_{(N-1) M}\right] \leq 1+\frac{p}{2 N} \int_{0}^{1} \mathbf{P}\left(B_{n}^{N-1} \leq x \mid \mathcal{F}_{(N-1) M}\right) x^{-p /(2 N)-1} \mathrm{~d} x \tag{22}
\end{equation*}
$$

Since $B_{n}^{N-1}$ is a polynomial of order $m=2(q-1)$, Carbery-Wright's inequality (18) yields:

$$
\begin{equation*}
\mathbf{P}\left(B_{n}^{N-1} \leq x \mid \mathcal{F}_{(N-1) M}\right) \leq \frac{c x^{1 /(2(q-1))}}{\left[\mathbf{E}\left(B_{n}^{N-1} \mid \mathcal{F}_{(N-1) M}\right)\right]^{1 /(2(q-1))}} . \tag{23}
\end{equation*}
$$

Step 4: Estimates for the conditional expectation. We now estimate the conditional expectation $\mathbf{E}\left[B_{n}^{N-1} \mid \mathcal{F}_{(N-1) M}\right]$. We have:

$$
\begin{align*}
\mathbf{E}\left[B_{n}^{N-1} \mid \mathcal{F}_{(N-1) M}\right] & =\frac{1}{n} \sum_{m=(N-1) M}^{N M-1} \mathbf{E}\left[\left(\sum_{k=0}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k}\right)^{2} \mid \mathcal{F}_{(N-1) M}\right]  \tag{24}\\
& \geq \frac{1}{n} \sum_{m=(N-1) M}^{N M-1} A_{m}
\end{align*}
$$

where we have set

$$
A_{m}=\operatorname{Var}\left(\sum_{k=(N-1) M}^{m} f^{\prime}\left(X_{k}\right) \psi_{m-k} \mid \mathcal{F}_{(N-1) M}\right)
$$

Furthermore, notice that

$$
f^{\prime}\left(X_{k}\right)=f^{\prime}\left(\sum_{\ell=-\infty}^{k} \psi_{k-i} w_{i}\right)=f^{\prime}\left(Y_{k}+Z_{k}\right)
$$

where $Y_{k}=\sum_{i=-\infty}^{(N-1) M-1} \psi_{k-i} w_{i}$ is $\mathcal{F}_{(N-1) M}$-measurable and $Z_{k}=\sum_{i=(N-1) M}^{k} \psi_{k-i} w_{i}$ is independent of $\mathcal{F}_{(N-1) M}$. Recalling that $f^{\prime}=\sum_{j=d}^{q} j a_{j} H_{j-1}$, we can thus resort to Lemmas 2.7
and 2.8. This gives:

$$
f^{\prime}\left(X_{k}\right)-\mathbf{E}\left[f^{\prime}\left(X_{k}\right) \mid \mathcal{F}_{(N-1) M}\right]=\sum_{j=d}^{q} \sum_{\ell=1}^{j-1} j a_{j}\binom{j-1}{\ell} \sigma_{Y_{k}}^{j-1-\ell} H_{j-1-\ell}\left(\widetilde{Y}_{k}\right) \sigma_{Z_{k}}^{\ell} H_{\ell}\left(\widetilde{Z}_{k}\right),
$$

where $\sigma_{Y_{k}}=\left[\operatorname{Var}\left(Y_{k}\right)\right]^{1 / 2}, \sigma_{Z_{k}}=\left[\operatorname{Var}\left(Z_{k}\right)\right]^{1 / 2}, \widetilde{Y}_{k}=Y_{k} / \sigma_{Y_{k}}$ and $\widetilde{Z}_{k}=Z_{k} / \sigma_{Z_{k}}$. Therefore,

$$
\begin{aligned}
A_{m} & =\mathbf{E}\left[\left(\sum_{k=(N-1) M}^{m} \sum_{j=d}^{q} \sum_{\ell=1}^{j-1} a_{j, \ell, k} H_{j-1-\ell}\left(\tilde{Y}_{k}\right) H_{\ell}\left(\widetilde{Z}_{k}\right) \psi_{m-k}\right)^{2} \mid \mathcal{F}_{(N-1) M}\right] \\
& =\mathbf{E}\left[\left(\sum_{\ell=1}^{q-1} \sum_{k=(N-1) M}^{m} \sum_{j=(\ell+1) \vee d}^{q} a_{j, \ell, k} H_{j-1-\ell}\left(\widetilde{Y}_{k}\right) H_{\ell}\left(\widetilde{Z}_{k}\right) \psi_{m-k}\right)^{2} \mid \mathcal{F}_{(N-1) M}\right],
\end{aligned}
$$

where we have set $a_{j, \ell, k}=j a_{j}\left({ }_{\ell}^{j-1}\right) \sigma_{Y_{k}}^{j-1-\ell} \sigma_{Z_{k}}^{\ell}$.
Recall that the random variables $\widetilde{Y}_{k}$ are $\mathcal{F}_{(N-1) M}$-measurable while the random variables $\widetilde{Z}_{k}$ are independent of $\mathcal{F}_{(N-1) M}$. By decorrelation properties of Hermite polynomials we thus get:

$$
A_{m}=\sum_{\ell=1}^{q-1} \mathbf{E}\left[\left(\sum_{k=(N-1) M}^{m} \sum_{j=(\ell+1) \vee d}^{q} a_{j, \ell, k} H_{j-1-\ell}\left(\widetilde{Y}_{k}\right) H_{\ell}\left(\widetilde{Z}_{k}\right) \psi_{m-k}\right)^{2} \mid \mathcal{F}_{(N-1) M}\right]
$$

and we trivially lower bound this quantity by taking the term corresponding to $\ell=q-1$. In this situation, the sum $\sum_{j=(\ell+1) \vee d}^{q}$ is reduced to the term corresponding to $j=q$, and since $a_{q, q-1, k}=q a_{q} \sigma_{Z_{k}}^{q-1}$ we obtain:

$$
\begin{aligned}
A_{m} & \geq \mathbf{E}\left[\left(\sum_{k=(N-1) M}^{m} q a_{q} \sigma_{Z_{k}}^{q-1} H_{q-1}\left(\widetilde{Z}_{k}\right) \psi_{m-k}\right)^{2} \mid \mathcal{F}_{(N-1) M}\right] \\
& =q^{2} a_{q}^{2} \mathbf{E}\left[\left(\sum_{k=(N-1) M}^{m} \sigma_{Z_{k}}^{q-1} H_{q-1}\left(\widetilde{Z}_{k}\right) \psi_{m-k}\right)^{2}\right] .
\end{aligned}
$$

We now invoke the identity $\mathbf{E}\left[H_{p}\left(\tilde{Z}_{k_{1}}\right) H_{p}\left(\tilde{Z}_{k_{2}}\right)\right]=p!\left(\mathbf{E}\left[\tilde{Z}_{k_{1}} \tilde{Z}_{k_{2}}\right]\right)^{p}$ in order to obtain

$$
A_{m} \geq q q!a_{q}^{2} \sum_{k_{1}, k_{2}=(N-1) M}^{m} \sigma_{Z_{k_{1}}}^{q-1} \sigma_{Z_{k_{2}}}^{q-1} \mathbf{E}\left[\widetilde{Z}_{k_{1}} \widetilde{Z}_{k_{2}}\right]^{q-1} \psi_{m-k_{1}} \psi_{m-k_{2}}
$$

Furthermore, similarly to (6), it is readily checked that:

$$
\mathbf{E}\left[\widetilde{Z}_{k_{1}} \widetilde{Z}_{k_{2}}\right]=\frac{1}{\sigma_{Z_{k_{1}}} \sigma_{Z_{k_{2}}}} \sum_{i=(N-1) M}^{k_{1} \wedge k_{2}} \psi_{k_{1}-i} \psi_{k_{2}-i}
$$

and thus

$$
\begin{aligned}
A_{m} & \geq q q!a_{q}^{2} \sum_{k_{1}, k_{2}=(N-1) M}^{m}\left(\sum_{i=(N-1) M}^{k_{1} \wedge k_{2}} \psi_{k_{1}-i} \psi_{k_{2}-i}\right)^{q-1} \psi_{m-k_{1}} \psi_{m-k_{2}} \\
& =q q!a_{q}^{2} \sum_{i_{1}, \ldots, i_{q-1}=(N-1) M}^{m} \sum_{k_{1}, k_{2}=\max \left(i_{1}, \ldots, i_{q-1}\right)}^{m} \psi_{m-k_{1}} \psi_{m-k_{2}} \prod_{j=1}^{q-1} \psi_{k_{1}-i_{j}} \psi_{k_{2}-i_{j}} \\
& =q q!a_{q}^{2} \sum_{i_{1}, \ldots, i_{q-1}=(N-1) M}^{m}\left(\sum_{k=\max \left(i_{1}, \ldots, i_{q-1}\right)}^{m} \psi_{m-k} \prod_{j=1}^{q-1} \psi_{k-i_{j}}\right)^{2}
\end{aligned}
$$

Here again, this sum of squares is trivially lower bounded by taking the term corresponding to $i_{1}=\cdots=i_{q-1}=m$, which yields:

$$
\begin{equation*}
A_{m} \geq c_{a, q, \psi} \quad \text { with } c_{a, q, \psi} \equiv q q!a_{q}^{2} \psi_{0}^{2 q}>0 \tag{25}
\end{equation*}
$$

Step 5: Conclusion. In the remainder of the proof the constants $c_{a, q, \psi, N}$ and so can change from line to line without further mention. Plugging relation (25) into (25) and recalling that $N$ is a given integer whose exact value will be fixed below, we get:

$$
\mathbf{E}\left[B_{n}^{N-1} \mid \mathcal{F}_{(N-1) M}\right] \geq \frac{M c_{a, q, \psi}}{n} \geq c_{a, q, \psi, N}>0
$$

as long as $N$ stays bounded. We then insert back this inequality into (22) and (23) in order to get:

$$
\mathbf{P}\left(B_{n}^{N-1} \leq x \mid \mathcal{F}_{(N-1) M}\right) \leq 1+\frac{p c_{a, q, \psi, N}}{2 N} \int_{0}^{1} x^{1 /(2(q-1))-p /(2 N)-1} \mathrm{~d} x=c_{a, q, \psi, N, p}<\infty
$$

where we have chosen $N$ such that $\frac{p}{2 N}<\frac{1}{2(q-1)}$. Iterating this bound into (21), we have thus obtained:

$$
\mathbf{E}\left[\left(B_{n}\right)^{-p / 2}\right] \leq c_{a, q, \psi, N, p}^{N}
$$

which is a finite quantity. Finally, recall from Step 1 that $\mathbf{E}\left[\left(B_{n}\right)^{-p / 2}\right] \geq \mathbf{E}\left[\left\|D V_{n}^{d, q}\right\|_{\mathfrak{H}}^{-p}\right]$, which finishes the proof.

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