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## To cite this version:

Ivan Nourdin, Giovanni Peccati. Stein's method on Wiener chaos. 2008. hal-00199024v4

## HAL Id: hal-00199024 https://hal.science/hal-00199024v4

Preprint submitted on 3 Feb 2008 (v4), last revised 10 May 2008 (v5)

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# Stein's method on Wiener chaos 

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Actual version: 3rd February 2008


#### Abstract

We combine Malliavin calculus with Stein's method, in order to derive explicit bounds in the Gaussian and Gamma approximations of random variables in a fixed Wiener chaos of a general Gaussian process. We also prove results concerning random variables admitting a possibly infinite Wiener chaotic decomposition. Our approach generalizes, refines and unifies the central and noncentral limit theorems for multiple Wiener-Itô integrals recently proved (in several papers, from 2005 to 2007) by Nourdin, Nualart, Ortiz-Latorre, Peccati and Tudor. We apply our techniques to prove Berry-Esséen bounds in the Breuer-Major CLT for subordinated functionals of fractional Brownian motion. By using the well-known Mehler's formula for Ornstein-Uhlenbeck semigroups, we also recover a result recently proved by Chatterjee, in the context of limit theorems for linear statistics of eigenvalues of random matrices.


Key words: Berry-Esséen bounds; Breuer-Major CLT; Fractional Brownian motion; Gamma approximation; Malliavin calculus; Multiple stochastic integrals; Normal approximation; Stein's method.

2000 Mathematics Subject Classification: 60F05; 60G15; 60H05; 60H07.

## 1 Introduction and overview

### 1.1 Motivations

Let $Z$ be a random variable whose law is absolutely continuous with respect to the Lebesgue measure (for instance, $Z$ is a standard Gaussian random variable or a Gamma random variable). Suppose that $\left\{Z_{n}: n \geqslant 1\right\}$ is a sequence of random variables converging in distribution towards $Z$, that is:

$$
\begin{equation*}
\text { for all } z \in \mathbb{R}, \quad P\left(Z_{n} \leqslant z\right) \longrightarrow P(Z \leqslant z) \quad \text { as } n \rightarrow \infty \tag{1.1}
\end{equation*}
$$

It is sometimes possible to associate an explicit uniform bound with the convergence (1.1), providing a global description of the error one makes when replacing $P(Z \leqslant z)$ by $P\left(Z_{n} \leqslant z\right)$ for a fixed $n \geqslant 1$. One of the most celebrated results in this direction is the following BerryEsséen Theorem (see e.g. Feller [16] for a proof), that we record here for future reference:

Theorem 1.1 (Berry-Esséen) Let $\left(U_{j}\right)_{j \geqslant 1}$ be a sequence of independent and identically distributed random variables, such that $E\left(\left|U_{j}\right|^{3}\right)=\rho<\infty, E\left(U_{j}\right)=0$ and $E\left(U_{j}^{2}\right)=\sigma^{2}$. Then,

[^0]by setting $Z_{n}=\frac{1}{\sigma \sqrt{n}} \sum_{j=1}^{n} U_{j}, n \geqslant 1$, one has that $Z_{n} \xrightarrow{\text { Law }} Z \sim \mathscr{N}(0,1)$, as $n \rightarrow \infty$, and moreover:
\[

$$
\begin{equation*}
\sup _{z \in \mathbb{R}}\left|P\left(Z_{n} \leqslant z\right)-P(Z \leqslant z)\right| \leqslant \frac{3 \rho}{\sigma^{3} \sqrt{n}} \tag{1.2}
\end{equation*}
$$

\]

The aim of this paper is to show that one can combine Malliavin calculus (see e.g. [34]) and Stein's method (see e.g. [10]), in order to obtain bounds analogous to (1.2), whenever the random variables $Z_{n}$ in (1.1) can be represented as functionals of a given Gaussian field. Our results are general, in the sense that (i) they do not rely on any specific assumption on the underlying Gaussian field, (ii) they do not require that the variables $Z_{n}$ have the specific form of partial sums, and (iii) they allow to deal (at least in the case of Gaussian approximations) with several different notions of distance between probability measures. As suggested by the title, a prominent role will be played by random variables belonging to a Wiener chaos of order $q(q \geqslant 2)$, that is, random variables having the form of a multiple stochastic WienerItô integral of order $q$ (see Section 2 below for precise definitions). It will be shown that our results provide substantial refinements of the central and non-central limit theorems for multiple stochastic integrals, recently proved in [32] and [36]. Among other applications and examples, we will provide explicit Berry-Esséen bounds in the Breuer-Major CLT (see [4]) for fields subordinated to a fractional Brownian motion.

Concerning point (iii), we shall note that, as a by-product of the flexibility of Stein's method, we will indeed establish bounds for Gaussian approximations related to a number of distances of the type

$$
\begin{equation*}
d_{\mathscr{H}}(X, Y)=\sup \{|E(h(X))-E(h(Y))|: h \in \mathscr{H}\}, \tag{1.3}
\end{equation*}
$$

where $\mathscr{H}$ is some suitable class of functions. For instance: by taking $\mathscr{H}=\left\{h:\|h\|_{L} \leqslant 1\right\}$, where $\|\cdot\|_{L}$ is the usual Lipschitz seminorm, one obtains the Wasserstein (or KantorovichWasserstein) distance; by taking $\mathscr{H}=\left\{h:\|h\|_{B L} \leqslant 1\right\}$, where $\|\cdot\|_{B L}=\|\cdot\|_{L}+\|\cdot\|_{\infty}$, one obtains the Fortet-Mourier (or bounded Wasserstein) distance; by taking $\mathscr{H}$ equal to the collection of all indicators $\mathbf{1}_{B}$ of Borel sets, one obtains the total variation distance; by taking $\mathscr{H}$ equal to the class of all indicators functions $\mathbf{1}_{(-\infty, z]}, z \in \mathbb{R}$, one has the Kolmogorov distance, which is the one taken into account in the Berry-Esséen bound (1.2). In what follows, we shall sometimes denote by $d_{\mathrm{W}}(.,),. d_{\mathrm{FM}}(.,),. d_{\mathrm{TV}}(.,$.$) and d_{\mathrm{Kol}}(.,$.$) , respectively,$ the Wasserstein, Fortet-Mourier, total variation and Kolmogorov distances. Observe that $d_{\mathrm{W}}(.,.) \geqslant d_{\mathrm{FM}}(.,$.$) and d_{\mathrm{TV}}(.,.) \geqslant d_{\mathrm{Kol}}(.,$.$) . Also, the topologies induced by d_{\mathrm{W}}, d_{\mathrm{TV}}$ and $d_{\text {Kol }}$ are stronger than the topology of convergence in distribution, while one can show that $d_{\mathrm{FM}}$ metrizes the convergence in distribution (see e.g. [15, Ch. 11] for these and further results involving distances on spaces of probability measures).

### 1.2 Stein's method

We shall now give a short account of Stein's method, which is basically a set of techniques allowing to evaluate distances of the type (1.3) by means of differential operators. This theory has been initiated by Stein in the path-breaking paper [46], and then further developed in the monograph [47]. The reader is referred to [10], [42] and [43] for detailed surveys of recent results and applications. The paper by Chatterjee [7] provides further insights into the existing literature. In what follows, we will apply Stein's method to two types of approximations,
namely Gaussian and (centered) Gamma. We shall denote by $\mathscr{N}(0,1)$ a standard Gaussian random variable. The centered Gamma random variables we are interested in have the form

$$
\begin{equation*}
F(\nu) \stackrel{\text { Law }}{=} 2 G(\nu / 2)-\nu, \quad \nu>0, \tag{1.4}
\end{equation*}
$$

where $G(\nu / 2)$ has a Gamma law with parameter $\nu / 2$. This means that $G(\nu / 2)$ is a (a.s. strictly positive) random variable with density $g(x)=\frac{x^{\frac{\nu}{2}-1} \mathrm{e}^{-x}}{\Gamma(\nu / 2)} \mathbf{1}_{(0, \infty)}(x)$, where $\Gamma$ is the usual Gamma function. We choose this parametrization in order to facilitate the connection with our previous paper [32] (observe in particular that, if $\nu \geqslant 1$ is an integer, then $F(\nu)$ has a centered $\chi^{2}$ distribution with $\nu$ degrees of freedom).

Standard Gaussian distribution. Let $Z \sim \mathscr{N}(0,1)$. Consider a real-valued function $h$ : $\mathbb{R} \rightarrow \mathbb{R}$ such that the expectation $E(h(Z))$ is well-defined. The Stein equation associated with $h$ and $Z$ is classically given by

$$
\begin{equation*}
h(x)-E(h(Z))=f^{\prime}(x)-x f(x), \quad x \in \mathbb{R} . \tag{1.5}
\end{equation*}
$$

A solution to (1.5) is a function $f$ which is Lebesgue a.e.-differentiable, and such that there exists a version of $f^{\prime}$ verifying (1.5) for every $x \in \mathbb{R}$. The following result is basically due to Stein [46, 47]. The proof of point (i) (whose content is usually referred as Stein's lemma) involves standard integration by parts (see e.g. [10, Lemma 2.1]). Point (ii) is proved e.g. in [10, Lemma 2.2]; point (iii) can be obtained by combining e.g. the arguments in [47, p. 25] and [8, Lemma 5.1]; a proof of point (iv) is contained in [47, Lemma 3, p. 25]; point (v) is proved in [7, Lemma 4.3].

Lemma 1.2 (i) Let $W$ be a random variable. Then, $W \stackrel{\text { Law }}{=} Z \sim \mathscr{N}(0,1)$ if, and only if,

$$
\begin{equation*}
E\left[f^{\prime}(W)-W f(W)\right]=0, \tag{1.6}
\end{equation*}
$$

for every continuous and piecewise continuously differentiable function $f$ verifying the relation $E\left|f^{\prime}(Z)\right|<\infty$.
(ii) If $h(x)=\mathbf{1}_{(-\infty, z]}(x)$, $z \in \mathbb{R}$, then (1.5) admits a solution $f$ which is bounded by $\sqrt{2 \pi} / 4$, piecewise continuously differentiable and such that $\left\|f^{\prime}\right\|_{\infty} \leqslant 1$.
(iii) If $h$ is bounded by $1 / 2$, then (1.5) admits a solution $f$ which is bounded by $\sqrt{\pi / 2}$, Lebesgue a.e. differentiable and such that $\left\|f^{\prime}\right\|_{\infty} \leqslant 2$.
(iv) If $h$ is bounded and absolutely continuous (then, in particular, Lebesgue-a.e. differentiable), then (1.5) has a solution $f$ which is bounded and twice differentiable, and such that $\|f\|_{\infty} \leqslant \sqrt{\pi / 2}\|h-E(h(Z))\|_{\infty},\left\|f^{\prime}\right\|_{\infty} \leqslant 2\|h-E(h(Z))\|_{\infty}$ and $\left\|f^{\prime \prime}\right\|_{\infty} \leqslant 2\left\|h^{\prime}\right\|_{\infty}$.
(v) If $h$ is absolutely continuous with bounded derivative, then (1.5) has a solution $f$ which is twice differentiable and such that $\left\|f^{\prime}\right\|_{\infty} \leqslant\left\|h^{\prime}\right\|_{\infty}$ and $\left\|f^{\prime \prime}\right\|_{\infty} \leqslant 2\left\|h^{\prime}\right\|_{\infty}$.

We also recall the relation:

$$
\begin{equation*}
2 d_{\mathrm{TV}}(X, Y)=\sup \left\{|E(u(X))-E(u(Y))|:\|u\|_{\infty} \leqslant 1\right\} . \tag{1.7}
\end{equation*}
$$

Note that point (ii) and (iii) (via (1.7)) imply the following bounds on the Kolmogorov and total variation distance between $Z$ and an arbitrary random variable $Y$ :

$$
\begin{align*}
& d_{\mathrm{Kol}}(Y, Z) \leqslant \sup _{f \in \mathscr{F}_{\mathrm{Kol}}}\left|E\left(f^{\prime}(Y)-Y f(Y)\right)\right|  \tag{1.8}\\
& d_{\mathrm{TV}}(Y, Z) \leqslant \sup _{f \in \mathscr{\mathscr { F }}_{\mathrm{TV}}}\left|E\left(f^{\prime}(Y)-Y f(Y)\right)\right| \tag{1.9}
\end{align*}
$$

where $\mathscr{F}_{\text {Kol }}$ and $\mathscr{F}_{\text {TV }}$ are, respectively, the class of piecewise continuously differentiable functions that are bounded by $\sqrt{2 \pi} / 4$ and such that their derivative is bounded by 1 , and the class of piecewise continuously differentiable functions that are bounded by $\sqrt{\pi / 2}$ and such that their derivative is bounded by 2 .

Analogously, by using (iv) and (v) along with the relation $\|h\|_{L}=\left\|h^{\prime}\right\|_{\infty}$, one obtains

$$
\begin{array}{r}
d_{\mathrm{FM}}(Y, Z) \leqslant \sup _{f \in \mathscr{F}_{\mathrm{FM}}}\left|E\left(f^{\prime}(Y)-Y f(Y)\right)\right|, \\
d_{\mathrm{W}}(Y, Z) \leqslant \sup _{f \in \mathscr{\mathscr { F }}_{\mathrm{W}}}\left|E\left(f^{\prime}(Y)-Y f(Y)\right)\right|, \tag{1.11}
\end{array}
$$

where: $\mathscr{F}_{\mathrm{FM}}$ is the class of twice differentiable functions that are bounded by $\sqrt{2 \pi}$, whose first derivative is bounded by 4 , and whose second derivative is bounded by 2 ; $\mathscr{F}_{\mathrm{W}}$ is the class of twice differentiable functions, whose first derivative is bounded by 1 and whose second derivative is bounded by 2 .

Centered Gamma distribution. Let $F(\nu)$ be as in (1.4). Consider a real-valued function $h: \mathbb{R} \rightarrow \mathbb{R}$ such that the expectation $E[h(F(\nu))]$ exists. The Stein equation associated with $h$ and $F(\nu)$ is:

$$
\begin{equation*}
h(x)-E[h(F(\nu))]=2(x+\nu) f^{\prime}(x)-x f(x), \quad x \in(-\nu,+\infty) . \tag{1.12}
\end{equation*}
$$

The following statement collects some slight variations around results proved by Stein [47], Diaconis and Zabell [14], Luk [24], Schoutens [45] and Pickett [40]. It is the "Gamma counterpart" of Lemma 1.2. The proof is detailed in Section 5 .

Lemma 1.3 (i) Let $W$ be a real-valued random variable (not necessarily with values in $(-\nu,+\infty)$ ) whose law admits a density with respect to the Lebesgue measure. Then, $W \stackrel{\text { Law }}{=} F(\nu)$ if, and only if,

$$
\begin{equation*}
E\left[2(W+\nu)_{+} f^{\prime}(W)-W f(W)\right]=0 \tag{1.13}
\end{equation*}
$$

where $a_{+}:=\max (a, 0)$, for every smooth function $f$ such that the mapping $x \mapsto 2(x+$ $\nu)_{+} f^{\prime}(x)-x f(x)$ is bounded.
(ii) If $|h(x)| \leqslant c \exp (a x)$ for every $x>-\nu$ and for some $c>0$ and $a<1 / 2$, and if $h$ is twice differentiable, then (1.12) has a solution $f$ which is bounded on $(-\nu,+\infty)$, differentiable and such that $\|f\|_{\infty} \leqslant 2\left\|h^{\prime}\right\|_{\infty}$ and $\left\|f^{\prime}\right\|_{\infty} \leqslant\left\|h^{\prime \prime}\right\|_{\infty}$.
(iii) Suppose that $\nu \geqslant 1$ is an integer. If $|h(x)| \leqslant c \exp (a x)$ for every $x>-\nu$ and for some $c>0$ and $a<1 / 2$, and if $h$ is twice differentiable with bounded derivatives, then (1.12) has a solution $f$ which is bounded on $(-\nu,+\infty)$, differentiable and such that $\|f\|_{\infty} \leqslant \sqrt{2 \pi / \nu}\|h\|_{\infty}$ and $\left\|f^{\prime}\right\|_{\infty} \leqslant \sqrt{2 \pi / \nu}\left\|h^{\prime}\right\|_{\infty}$.

Now define

$$
\begin{align*}
\mathscr{H}_{1} & =\left\{h \in \mathscr{C}_{b}^{2}:\|h\|_{\infty} \leqslant 1,\left\|h^{\prime}\right\|_{\infty} \leqslant 1,\left\|h^{\prime \prime}\right\|_{\infty} \leqslant 1\right\}  \tag{1.14}\\
\mathscr{H}_{2} & =\left\{h \in \mathscr{C}_{b}^{2}:\|h\|_{\infty} \leqslant 1,\left\|h^{\prime}\right\|_{\infty} \leqslant 1\right\}  \tag{1.15}\\
\mathscr{H}_{1, \nu} & =\mathscr{H}_{1} \cap \mathscr{C}_{b}^{2}(\nu)  \tag{1.16}\\
\mathscr{H}_{2, \nu} & =\mathscr{H}_{2} \cap \mathscr{C}_{b}^{2}(\nu) \tag{1.17}
\end{align*}
$$

where $\mathscr{C}_{b}^{2}$ denotes the class of twice differentiable functions with support in $\mathbb{R}$ and with bounded derivatives, and $\mathscr{C}_{b}^{2}(\nu)$ denotes the subset of $\mathscr{C}_{b}^{2}$ composed of functions with support in $(-\nu,+\infty)$. Note that point (ii) in the previous statement implies that, by adopting the notation (1.3) and for every $\nu>0$ and every real random variable $Y$ (not necessarily with support in $(-\nu,+\infty))$,

$$
\begin{equation*}
d_{\mathscr{H}_{1, \nu}}(Y, F(\nu)) \leqslant \sup _{f \in \mathscr{F}_{1, \nu}}\left|E\left[2(Y+\nu) f^{\prime}(Y)-Y f(Y)\right]\right| \tag{1.18}
\end{equation*}
$$

where $\mathscr{F}_{1, \nu}$ is the class of differentiable functions with support in $(-\nu,+\infty)$, bounded by 2 and whose derivatives are bounded by 1. Analogously, point (iii) implies that, for every integer $\nu \geqslant 1$,

$$
\begin{equation*}
d_{\mathscr{H}_{2, \nu}}(Y, F(\nu)) \leqslant \sup _{f \in \mathscr{F}_{2, \nu}}\left|E\left[2(Y+\nu) f^{\prime}(Y)-Y f(Y)\right]\right| \tag{1.19}
\end{equation*}
$$

where $\mathscr{F}_{2, \nu}$ is the class of differentiable functions with support in $(-\nu,+\infty)$, bounded by $\sqrt{2 \pi / \nu}$ and whose derivatives are also bounded by $\sqrt{2 \pi / \nu}$. A little inspection shows that the following estimates also hold: for every $\nu>0$ and every random variable $Y$,

$$
\begin{equation*}
d_{\mathscr{H}_{1}}(Y, F(\nu)) \leqslant \sup _{f \in \mathscr{F}_{1}}\left|E\left[2(Y+\nu)_{+} f^{\prime}(Y)-Y f(Y)\right]\right| \tag{1.20}
\end{equation*}
$$

where $\mathscr{F}_{1}$ is the class of functions (defined on $\mathbb{R}$ ) that are continuous and differentiable on $\mathbb{R} \backslash\{\nu\}$, bounded by $\max \{2,2 / \nu\}$, and whose derivatives are bounded by $\max \left\{1,1 / \nu+2 / \nu^{2}\right\}$. Analogously, for every integer $\nu \geqslant 1$,

$$
\begin{equation*}
d_{\mathscr{H}_{2}}(Y, F(\nu)) \leqslant \sup _{f \in \mathscr{F}_{2}}\left|E\left[2(Y+\nu)_{+} f^{\prime}(Y)-Y f(Y)\right]\right|, \tag{1.21}
\end{equation*}
$$

where $\mathscr{F}_{2}$ is the class of functions (on $\mathbb{R}$ ) that are continuous and differentiable on $\mathbb{R} \backslash\{\nu\}$, bounded by $\max \{\sqrt{2 \pi / \nu}, 2 / \nu\}$, and whose derivatives are bounded by $\max \{\sqrt{2 \pi / \nu}, 1 / \nu+$ $\left.2 / \nu^{2}\right\}$.

Now, the crucial issue is how to estimate the right-hand side of (1.8)-(1.11) and (1.18)(1.21) for a given choice of $Y$. Since Stein's initial contribution [46], an impressive panoply of techniques has been developed in this direction (see again [9] or [42] for a survey; here, we shall quote e.g.: exchangeable pairs [47], diffusion generators [3, 18], size-bias transforms [19], zerobias transforms [20], local dependency graphs [9] and graphical-geometric rules [7]). Starting from the next section, we will show that, when working within the framework of functionals of Gaussian fields, one can very effectively estimate expressions such as (1.8)-(1.11), (1.18) and (1.19) by using techniques of Malliavin calculus. Interestingly, a central role is played by an infinite dimensional version of the same integration by parts formula that is at the very heart of Stein's characterization of the Gaussian distribution.

### 1.3 The basic approach (with some examples)

Let $\mathfrak{H}$ be a real separable Hilbert space and, for $q \geqslant 1$, let $\mathfrak{H}^{\otimes q}$ (resp. $\mathfrak{H}^{\odot q}$ ) be the $q$ th tensor product (resp. $q$ th symmetric tensor product) of $\mathfrak{H}$. We write

$$
\begin{equation*}
X=\{X(h): h \in \mathfrak{H}\} \tag{1.22}
\end{equation*}
$$

to indicate a centered isonormal Gaussian process on $\mathfrak{H}$. For every $q \geqslant 1$, we denote by $I_{q}$ the isometry between $\mathfrak{H}^{\odot q}$ (equipped with the norm $\sqrt{q!}\|\cdot\|_{\mathfrak{S}^{\otimes q}}$ ) and the $q$ th Wiener chaos of $X$. Note that, if $\mathfrak{H}$ is a $\sigma$-finite measure space with no atoms, then each random variable $I_{q}(h), h \in \mathfrak{H}^{\odot q}$, has the form of a multiple Wiener-Itô integral of order $q$. We denote by $\mathrm{L}^{2}(X)=\mathrm{L}^{2}(\Omega, \sigma(X), P)$ the space of square integrable functionals of $X$, and by $\mathbb{D}^{1,2}$ the domain of the Malliavin derivative operator $D$ (see the forthcoming Section 2 for precise definitions). Recall that, for every $F \in \mathbb{D}^{1,2}, D F$ is a random element with values in $\mathfrak{H}$.

We start by observing that, thanks to (1.6), for every $h \in \mathfrak{H}$ such that $\|h\|_{\mathfrak{H}}=1$ and for every smooth function $f$, we have $E[X(h) f(X(h))]=E\left[f^{\prime}(X(h))\right]$. Our point is that this last relation is a very particular case of the following corollary of the celebrated integration by parts formula of Malliavin calculus: for every $Y \in \mathbb{D}^{1,2}$,

$$
\begin{equation*}
E[Y f(Y)]=E\left[\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{s}} f^{\prime}(Y)\right], \tag{1.23}
\end{equation*}
$$

where the linear operator $L^{-1}$, from $L^{2}(X)$ into $\mathbb{D}^{2,2} \subset \mathbb{D}^{1,2}$ (the space $\mathbb{D}^{2,2}$ is defined in Section 2), is the inverse of the generator of the Ornstein-Uhlebeck semigroup, noted $L$. The reader is referred to Section 2 and Section 3 for definitions and for a full discussion of this point; here, we shall note that $L$ is an infinite-dimensional version of the generator associated with Ornstein-Uhlenbeck diffusions (see [34, Section 1.4] for a proof of this fact), an object which is also crucial in the Barbour-Götze "generator approach" to Stein's method [3, 18].

It follows that, for every $Y \in \mathbb{D}^{1,2}$, the expressions appearing on the RHS of (1.8)-(1.11) (or (1.18)-(1.21)) can be assessed by first replacing $Y f(Y)$ with $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{f}} f^{\prime}(Y)$ inside the expectation, and then by evaluating the $\mathrm{L}^{2}$ distance between 1 (resp. $2 Y+2 \nu$ ) and the inner product $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{j}}$. In general, this computations are carried out by first resorting to the representation of $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{H}}$ as a (possibly infinite) series of multiple stochastic integrals. We will see that, when $Y=I_{q}(g)$, for $q \geqslant 2$ and some $g \in \mathfrak{H}^{\odot q}$, then $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{H}}=q^{-1}\|D Y\|_{\mathfrak{H}}^{2}$. In particular, by using this last relation one can deduce bounds involving quantities that are intimately related to the central and non-central limit theorems recently proved in [36], [35] and [32].

Remark 1.4 1. The crucial equality $E\left[I_{q}(g) f\left(I_{q}(g)\right)\right]=E\left[q^{-1}\left\|D I_{q}(g)\right\|_{\mathfrak{S}}^{2} f^{\prime}\left(I_{q}(g)\right)\right]$, in the case where $f$ is a complex exponential, has been first used in [35], in order to give refinements (as well as alternate proofs) of the main CLTs in [36] and [38]. The same relation has been later applied in [32], where a characterization of non-central limit theorems for multiple integrals is provided. Note that neither [32] nor [35] are concerned with Stein's method or, more generally, with bounds on distances between probability measures.
2. We will see that formula (1.23) contains as a special case a result recently proved by Chatterjee [8, Lemma 5.3], in the context of limit theorems for linear statistics of eigenvalues of random matrices. The connection between the two results can be established by means of the well-known Mehler's formula (see e.g. [26, Section 8.5, Ch. I] or [34, Section 1.4]), providing a mixture-type representation of the infinite-dimensional Ornstein-Uhlenbeck semigroup. See Remarks 3.6 and 3.12 below for a precise discussion of this point. See e.g. [29] for a detailed presentation of the infinite-dimensional Ornstein-Uhlebeck semigroup.
3. We stress that the random variable $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{5}}$ appearing in (1.23) is in general not measurable with respect to $\sigma(Y)$. For instance, if $X$ is taken to be the Gaussian
space generated by a standard Brownian motion $\left\{W_{t}: t \geqslant 0\right\}$ and $Y=I_{2}(h)$ with $h \in L_{s}^{2}\left([0,1]^{2}\right)$, then $D_{t} Y=2 \int_{0}^{1} h(u, t) d W_{u}, t \in[0,1]$, and

$$
\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathrm{L}^{2}([0,1])}=2 I_{2}\left(h \otimes_{1} h\right)+2\|h\|_{\mathrm{L}^{2}\left([0,1]^{2}\right)}^{2}
$$

which is, in general, not measurable with respect to $\sigma(Y)$ (the symbol $h \otimes_{1} h$ indicates a contraction kernel that will be defined in Section 2).
4. Note that (1.23) also implies the relation

$$
\begin{equation*}
E[Y f(Y)]=E\left[\tau(Y) f^{\prime}(Y)\right] \tag{1.24}
\end{equation*}
$$

where $\tau(Y)=E\left[\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{H}} \mid Y\right]$. Some general results for the existence of a realvalued function $\tau$ satisfying (1.24) are contained e.g. in [47, Lecture VI] and [5]. Note that, in general, it is very hard to find an analytic expression for $\tau(Y)$, especially when $Y$ is a random variable with a very complex structure, such as e.g. a multiple WienerItô integral. On the other hand, we will see that, in many cases, the random variable $\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{H}}$ is remarkably tractable and explicit. See Remark 3.10 below for a connection with Goldstein and Reinert's zero bias transform [19].
5. The reader is referred to [41] for applications of integration by parts techniques to the Stein-type estimation of drifts of Gaussian processes. See [21] for a Stein characterization of Brownian motions on manifolds by means of integration by parts formulae. See [12] for a connection between Stein's method and algebras of operators on configuration spaces.

Before proceeding to a formal discussion, and in order to motivate the reader, we shall provide two examples of the kind of results that we will obtain in the subsequent sections. The first statement involves double Wiener-Itô integrals, that is, random variables living in the second chaos of $X$. The proof is given in Section 5 .

Theorem 1.5 Let $\left(Z_{n}\right)_{n \geqslant 1}$ be a sequence belonging to the second Wiener chaos of $X$.
 $n \rightarrow \infty$. Moreover, we have:

$$
d_{\mathrm{TV}}\left(Z_{n}, Z\right) \leqslant 2 \sqrt{\frac{1}{6}\left|E\left(Z_{n}^{4}\right)-3\right|+\frac{3+E\left(Z_{n}^{2}\right)}{2}\left|E\left(Z_{n}^{2}\right)-1\right|}
$$

2. Fix $\nu>0$ and assume that $E\left(Z_{n}^{2}\right) \rightarrow 2 \nu$ and $E\left(Z_{n}^{4}\right)-12 E\left(Z_{n}^{3}\right) \rightarrow 12 \nu^{2}-48 \nu$ as $n \rightarrow \infty$. Then, as $n \rightarrow \infty, Z_{n} \xrightarrow{\text { Law }} F(\nu)$, where $F(\nu)$ has a centered Gamma distribution of parameter ע. Moreover, we have:

$$
\begin{aligned}
& d_{\mathscr{H}_{1}}\left(Z_{n}, F(\nu)\right) \\
\leqslant & K_{1} \sqrt{\frac{1}{6}\left|E\left(Z_{n}^{4}\right)-12 E\left(Z_{n}^{3}\right)-12 \nu^{2}+48 \nu\right|+\frac{\left|8-6 \nu+E\left(Z_{n}^{2}\right)\right|}{2}\left|E\left(Z_{n}^{2}\right)-2 \nu\right|},
\end{aligned}
$$

where $K_{1}:=\max \left\{1,1 / \nu, 2 / \nu^{2}\right\}$ and $\mathscr{H}_{1}$ is defined by (1.14).
For instance, when applied to the case where $X$ is the isonormal process generated by a fractional Brownian motion, the first point of Theorem 1.5 can be used to derive the following bound for the Kolmogorov distance in the Breuer-Major CLT associated with quadratic transformations:

Theorem 1.6 Let $B$ be a fractional Brownian motion with Hurst index $H \in(0,3 / 4)$. We set

$$
\sigma_{H}^{2}=\frac{1}{2} \sum_{t \in \mathbb{Z}}\left(|t+1|^{2 H}+|t-1|^{2 H}-2|t|^{2 H}\right)^{2}<\infty
$$

and

$$
Z_{n}=\frac{1}{\sigma_{H} \sqrt{n}} \sum_{k=0}^{n-1}\left(n^{2 H}\left(B_{(k+1) / n}-B_{k / n}\right)^{2}-1\right), \quad n \geqslant 1 .
$$

Then, as $n \rightarrow \infty, Z_{n} \xrightarrow{\text { Law }} Z \sim \mathscr{N}(0,1)$. Moreover, there exists a constant $c_{H}$ (depending uniquely on $H$ ) such that, for any $n \geqslant 1$ :

$$
\begin{equation*}
d_{\mathrm{Kol}}\left(Z_{n}, Z\right) \leqslant \frac{c_{H}}{n^{\frac{1}{2} \wedge\left(\frac{3}{2}-2 H\right)}} \tag{1.25}
\end{equation*}
$$

Note that both Theorem 1.5 and 1.6 will be significantly generalized in Section 3 and Section 4.

Remark 1.7 1. When $H=1 / 2$, then $B$ is a standard Brownian motion (and therefore has independent increments), and we recover from the previous result the rate $n^{-1 / 2}$, that could be also obtained by applying the Berry-Esséen Theorem 1.1. This rate is still valid for $H<1 / 2$. But, for $H>1 / 2$, the rate in the Breuer-Major CLT becomes $n^{2 H-\frac{3}{2}}$.
2. To the authors knowledge, Theorem 1.6 and its generalizations are the first Berry-Esséen bounds ever established for the Breuer-Major CLT.
3. To keep the length of this paper within limits, we do not derive the explicit expression of some of the constants (such as the quantity $c_{H}$ in formula (1.25)) composing our bounds. As will become clear later on, the exact value of these quantities can be deduced by a careful bookkeeping of the bounding constants appearing at the different stages of the proofs.

### 1.4 Plan

The rest of the paper is organized as follows. In Section 2 we recall some notion of Malliavin calculus; in Section 3 we state and discuss our main bounds in Stein-type estimates for functionals of Gaussian fields; Section 4 contains an application to the Breuer-Major CLT. Proofs and further refinements are collected in Section 5.

## 2 Elements of Malliavin calculus

The reader is referred to [23] or [34] for any unexplained notion discussed in this section. As in (1.22), we denote by $X=\{X(h): h \in \mathfrak{H}\}$ an isonormal Gaussian process over $\mathfrak{H}$. By definition, $X$ is a centered Gaussian family indexed by the elements of $\mathfrak{H}$ and such that, for every $h, g \in \mathfrak{H}$,

$$
\begin{equation*}
E[X(h) X(g)]=\langle h, g\rangle_{\mathfrak{H}} . \tag{2.26}
\end{equation*}
$$

As before, we use the notation $\mathrm{L}^{2}(X)=\mathrm{L}^{2}(\Omega, \sigma(X), P)$. It is well-known (see again [34, Ch. 1] or [23]) that any random variable $F$ belonging to $\mathrm{L}^{2}(X)$ admits the following chaotic expansion:

$$
\begin{equation*}
F=\sum_{q=0}^{\infty} I_{q}\left(f_{q}\right) \tag{2.27}
\end{equation*}
$$

where $I_{0}\left(f_{0}\right):=E[F]$, the series converges in $\mathrm{L}^{2}$ and the kernels $f_{q} \in \mathfrak{H}^{\odot q}, q \geqslant 1$, are uniquely determined by $F$. As already discussed, in the particular case where $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$, where $(A, \mathscr{A})$ is a measurable space and $\mu$ is a $\sigma$-finite and non-atomic measure, one has that $\mathfrak{H}^{\odot q}=\mathrm{L}_{s}^{2}\left(A^{q}, \mathscr{A}^{\otimes q}, \mu^{\otimes q}\right)$ is the space of symmetric and square integrable functions on $A^{q}$. Moreover, for every $f \in \mathfrak{H}^{\odot q}, I_{q}(f)$ coincides with the multiple Wiener-Itô integral (of order $q)$ of $f$ with respect to $X$ (see [34, Ch. 1]). Observe that a random variable of the type $I_{q}(f)$, $f \in \mathfrak{H}^{\odot q}$, has finite moments of all orders (see e.g. [23, Ch. VI]). See again [34, Ch. 1] or [44] for a connection between multiple Wiener-Itô and Hermite polynomials. For every $q \geqslant 0$, we write $J_{q}$ to indicate the orthogonal projection operator on the $q$ th Wiener chaos associated with $X$, so that, if $F \in \mathrm{~L}^{2}(\Omega, \mathscr{F}, P)$ is as in (2.27), then $J_{q} F=I_{q}\left(f_{q}\right)$ for every $q \geqslant 0$.

Let $\left\{e_{k}, k \geq 1\right\}$ be a complete orthonormal system in $\mathfrak{H}$. Given $f \in \mathfrak{H}^{\odot p}$ and $g \in \mathfrak{H}^{\odot q}$, for every $r=0, \ldots, p \wedge q$, the $r$ th contraction of $f$ and $g$ is the element of $\mathfrak{H}^{\otimes(p+q-2 r)}$ defined as

$$
\begin{equation*}
f \otimes_{r} g=\sum_{i_{1}, \ldots, i_{r}=1}^{\infty}\left\langle f, e_{i_{1}} \otimes \ldots \otimes e_{i_{r}}\right\rangle_{\mathfrak{H}^{\otimes r}} \otimes\left\langle g, e_{i_{1}} \otimes \ldots \otimes e_{i_{r}}\right\rangle_{\mathfrak{H} \otimes r} \tag{2.28}
\end{equation*}
$$

Note that, in the particular case where $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$ (with $\mu$ non-atomic), one has that

$$
f \otimes_{r} g=\int_{A^{r}} f\left(t_{1}, \ldots, t_{p-r}, s_{1}, \ldots, s_{r}\right) g\left(t_{p-r+1}, \ldots, t_{p+q-2 r}, s_{1}, \ldots, s_{r}\right) d \mu\left(s_{1}\right) \ldots d \mu\left(s_{r}\right)
$$

Moreover, $f \otimes_{0} g=f \otimes g$ equals the tensor product of $f$ and $g$ while, for $p=q, f \otimes_{p} g=$ $\langle f, g\rangle_{\mathfrak{H}^{\otimes p}}$. Note that, in general (and except for trivial cases), the contraction $f \otimes_{r} g$ is not a symmetric element of $\mathfrak{H}^{\otimes(p+q-2 r)}$. The canonical symmetrization of $f \otimes_{r} g$ is written $f \widetilde{\otimes}_{r} g$. We also have the useful multiplication formula: if $f \in \mathfrak{H}^{\odot p}$ and $g \in \mathfrak{H}^{\odot q}$, then

$$
\begin{equation*}
I_{p}(f) I_{q}(g)=\sum_{r=0}^{p \wedge q} r!\binom{p}{r}\binom{q}{r} I_{p+q-2 r}\left(f \widetilde{\otimes}_{r} g\right) \tag{2.29}
\end{equation*}
$$

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

$$
F=g\left(X\left(\phi_{1}\right), \ldots, X\left(\phi_{n}\right)\right)
$$

where $n \geqslant 1, g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a smooth function with compact support and $\phi_{i} \in \mathfrak{H}$. The Malliavin derivative of $F$ with respect to $X$ is the element of $\mathrm{L}^{2}(\Omega, \mathfrak{H})$ defined as

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

In particular, $D X(h)=h$ for every $h \in \mathfrak{H}$. By iteration, one can define the $m$ th derivative $D^{m} F$ (which is an element of $\mathbb{L}^{2}\left(\Omega, \mathfrak{H}^{\otimes m}\right)$ ) for every $m \geqslant 2$. As usual, for $m \geqslant 1, \mathbb{D}^{m, 2}$ denotes the closure of $\mathscr{S}$ with respect to the norm $\|\cdot\|_{m, 2}$, defined by the relation

$$
\|F\|_{m, 2}^{2}=E\left[F^{2}\right]+\sum_{i=1}^{m} E\left[\left\|D^{i} F\right\|_{\mathfrak{H}^{\otimes i}}^{2}\right]
$$

Note that, if $F \neq 0$ and $F$ is aequal to a finite sum of multiple Wiener-Itô integrals, then $F \in \mathbb{D}^{m, 2}$ for every $m \geqslant 1$ and the law of $F$ admits a density with respect to Lebesgue measure. The Malliavin derivative $D$ verifies the following chain rule: if $\varphi: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is in $\mathscr{C}_{b}^{1}$ (that is, the collection of bounded 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

$$
D \varphi\left(F_{1}, \ldots, F_{n}\right)=\sum_{i=1}^{n} \frac{\partial \varphi}{\partial x_{i}}\left(F_{1}, \ldots, F_{n}\right) D F_{i} .
$$

Observe that the previous formula still holds when $\varphi$ is a Lipschitz function and the law of $\left(F_{1}, \ldots, F_{n}\right)$ has a density with respect to the Lebesgue measure on $\mathbb{R}^{n}$ (see e.g. Proposition 1.2.3 in [34]). We denote by $\delta$ the adjoint of the operator $D$, also called the divergence operator. A random element $u \in \mathrm{~L}^{2}(\Omega, \mathfrak{H})$ belongs to the domain of $\delta$, noted Dom $\delta$, if, and only if, it verifies

$$
\left|E\langle D F, u\rangle_{\mathfrak{H}}\right| \leqslant c_{u}\|F\|_{\mathrm{L}^{2}} \quad \text { for any } F \in \mathscr{S},
$$

where $c_{u}$ is a constant depending uniquely on $u$. If $u \in \operatorname{Dom} \delta$, then the random variable $\delta(u)$ is defined by the duality relationship (customarily called "integration by parts formula"):

$$
\begin{equation*}
E(F \delta(u))=E\langle D F, u\rangle_{\mathfrak{H}}, \tag{2.30}
\end{equation*}
$$

which holds for every $F \in \mathbb{D}^{1,2}$. One sometimes needs the following property: for every $F \in \mathbb{D}^{1,2}$ and every $u \in \operatorname{Dom} \delta$ such that $F u$ and $F \delta(u)+\langle D F, u\rangle_{\mathfrak{H}}$ are square integrable, one has that $F u \in \operatorname{Dom} \delta$ and

$$
\begin{equation*}
\delta(F u)=F \delta(u)-\langle D F, u\rangle_{\mathfrak{5}} . \tag{2.31}
\end{equation*}
$$

The operator $L$, acting on square integrable random variables of the type (2.27), is defined through the projection operators $\left\{J_{q}\right\}_{q \geqslant 0}$ as $L=\sum_{q=0}^{\infty}-q J_{q}$, and is called the infinitesimal generator of the Ornstein-Uhlenbeck semigroup. It verifies the following crucial property: a random variable $F$ is an element of $\operatorname{Dom} L\left(=\mathbb{D}^{2,2}\right)$ if, and only if, $F \in \operatorname{Dom} \delta D$ (i.e. $F \in \mathbb{D}^{1,2}$ and $D F \in \operatorname{Dom} \delta)$, and in this case: $\delta D F=-L F$. Note that a random variable $F$ as in (2.27) is in $\mathbb{D}^{1,2}$ (resp. $\mathbb{D}^{2,2}$ ) if, and only if,

$$
\left.\sum_{q=1}^{\infty} q\left\|f_{q}\right\|_{\mathfrak{S}^{\otimes q}}^{2}<\infty \quad \text { (resp. } \sum_{q=1}^{\infty} q^{2}\left\|f_{q}\right\|_{\mathfrak{S}^{\otimes q}}^{2}<\infty\right)
$$

and also $E\left[\|D F\|_{\mathfrak{H}}^{2}\right]=\sum_{q \geqslant 1} q\left\|f_{q}\right\|_{\mathfrak{H}^{\otimes q}}^{2}$. If $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$ (with $\mu$ non-atomic), then the derivative of a random variable $F$ as in (2.27) can be identified with the element of $L^{2}(A \times \Omega)$ given by

$$
\begin{equation*}
D_{a} F=\sum_{q=1}^{\infty} q I_{q-1}\left(f_{q}(\cdot, a)\right), \quad a \in A \tag{2.32}
\end{equation*}
$$

We also define the operator $L^{-1}$, which is the inverse of $L$, as follows: for every $F \in \mathrm{~L}^{2}(X)$ with zero mean, we set $L^{-1} F=\sum_{q \geqslant 1}-\frac{1}{q} J_{q}(F)$. Note that $L^{-1}$ is an operator with values in $\mathbb{D}^{2,2}$. The following Lemma contains two statements: the first one (formula (2.33)) is an immediate consequence of the definition of $L$ and of the relation $\delta D=-L$, whereas the second (formula (2.34)) corresponds to Lemma 2.1 in [32].

Lemma 2.1 Fix an integer $q \geqslant 2$ and set $F=I_{q}(f)$, with $f \in \mathfrak{H}^{\odot q}$. Then,

$$
\begin{equation*}
\delta D F=q F \tag{2.33}
\end{equation*}
$$

Moreover, for every integer $s \geqslant 0$,

$$
\begin{equation*}
E\left(F^{s}\|D F\|_{\mathfrak{H}}^{2}\right)=\frac{q}{s+1} E\left(F^{s+2}\right) \tag{2.34}
\end{equation*}
$$

## 3 Stein's method and integration by parts on Wiener space

### 3.1 Gaussian approximations

Our first result provides explicit bounds for the normal approximation of random variables that are Malliavin derivable. Although its proof is quite easy to obtain, the following statement will be central for the rest of the paper.

Theorem 3.1 Let $Z \sim \mathscr{N}(0,1)$, and let $F \in \mathbb{D}^{1,2}$ be such that $E(F)=0$. Then, the following bounds are in order:

$$
\begin{align*}
d_{\mathrm{W}}(F, Z) & \leqslant E\left[\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2}  \tag{3.35}\\
d_{\mathrm{FM}}(F, Z) & \leqslant 4 E\left[\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2} \tag{3.36}
\end{align*}
$$

If, in addition, the law of $F$ is absolutely continuous with respect to Lebesgue measure, one has that

$$
\begin{align*}
d_{\mathrm{Kol}}(F, Z) & \leqslant E\left[\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2}  \tag{3.37}\\
d_{\mathrm{TV}}(F, Z) & \leqslant 2 E\left[\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2} \tag{3.38}
\end{align*}
$$

Proof. Start by observing that one can write $F=L L^{-1} F=-\delta D L^{-1} F$. Now let $f$ be a real differentiable function. By using the integration by parts formula and the fact that $D f(F)=f^{\prime}(F) D F$ (note that, for this formula to hold when $f$ is only a.e. differentiable, one needs $F$ to have an absolutely continuous law, see Proposition 1.2.3 in [34]), we deduce

$$
E(F f(F))=E\left[f^{\prime}(F)\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right]
$$

It follows that $E\left[f^{\prime}(F)-F f(F)\right]=E\left(f^{\prime}(F)\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)\right)$ so that relations (3.35)(3.38) can be deduced from (1.8)-(1.11) and the Cauchy-Schwarz inequality.

We shall now prove that the bounds appearing in the statement of Theorem 3.1 can be explicitly computed, whenever $F$ belongs to a fixed Wiener chaos.

Proposition 3.2 Let $q \geqslant 2$ be an integer, and let $F=I_{q}(f)$, where $f \in \mathfrak{H}^{\odot q}$. Then, $\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}=q^{-1}\|D F\|_{\mathfrak{H}}^{2}$, and

$$
\begin{align*}
& E\left[\left(1-\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}\right)^{2}\right]=E\left[\left(1-q^{-1}\|D F\|_{\mathfrak{H}}^{2}\right)^{2}\right]  \tag{3.39}\\
& =\left(1-q!\|f\|_{\mathfrak{H}^{\otimes q}}^{2}\right)^{2}+q^{2} \sum_{r=1}^{q-1}(2 q-2 r)!(r-1)!^{2}\binom{q-1}{r-1}^{4}\left\|f \widetilde{\otimes}_{r} f\right\|_{\mathfrak{H}^{\otimes 2(q-r)}}^{2}  \tag{3.40}\\
& \leqslant\left(1-q!\|f\|_{\mathfrak{H}^{\otimes q}}^{2}\right)^{2}+q^{2} \sum_{r=1}^{q-1}(2 q-2 r)!(r-1)!^{2}\binom{q-1}{r-1}^{4}\left\|f \otimes_{r} f\right\|_{\mathfrak{H}^{\otimes 2(q-r)}}^{2} . \tag{3.41}
\end{align*}
$$

Proof. The equality $\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}=q^{-1}\|D F\|_{\mathfrak{H}}^{2}$ is an immediate consequence of the relation $L^{-1} I_{q}(f)=-q^{-1} I_{q}(f)$. From the multiplication formulae between multiple stochastic integrals, see (2.29), one deduces that

$$
\begin{equation*}
\left\|D\left[I_{q}(f)\right]\right\|_{\mathfrak{H}}^{2}=q q!\|f\|_{\mathfrak{H} \otimes q}^{2}+q^{2} \sum_{r=1}^{q-1}(r-1)!\binom{q-1}{r-1}^{2} I_{2(q-r)}\left(f \widetilde{\otimes}_{r} f\right) \tag{3.42}
\end{equation*}
$$

(see also [35, Lemma 2]). We therefore obtain (3.40) by using the orthogonality and isometric properties of multiple stochastic integrals. The inequality in (3.41) is just a consequence of the relation $\left\|f \widetilde{\otimes}_{r} f\right\|_{\mathfrak{H}^{\otimes 2(q-r)}} \leqslant\left\|f \otimes_{r} f\right\|_{\mathfrak{H}^{\otimes 2(q-r)}}$.

The previous result should be compared with the forthcoming Theorem 3.3, where we collect the main findings of [35] and [36]. In particular, the combination of Proposition 3.2 and Theorem 3.3 shows that, for every (normalized) sequence $\left\{F_{n}: n \geqslant 1\right\}$ living in a fixed Wiener chaos, the bounds given in (3.35)-(3.36) are "optimal" with respect to the convergence in distribution towards $Z \sim \mathscr{N}(0,1)$, in the sense that these bounds converge to zero if, and only if, $F_{n}$ converges in distribution to $Z$.

Theorem $3.3([35,36])$ Fix $q \geqslant 2$, and consider a sequence $\left\{F_{n}: n \geqslant 1\right\}$ such that $F_{n}=$ $I_{q}\left(f_{n}\right), n \geqslant 1$, where $f_{n} \in \mathfrak{H}^{\odot q}$. Assume moreover that $E\left[F_{n}^{2}\right]=q!\left\|f_{n}\right\|_{\mathfrak{H}^{\otimes q}}^{2} \rightarrow 1$. Then, the following four conditions are equivalent, as $n \rightarrow \infty$ :
(i) $F_{n}$ converges in distribution to $Z \sim \mathscr{N}(0,1)$;
(ii) $E\left[F_{n}^{4}\right] \rightarrow 3$;
(iii) for every $r=1, \ldots, q-1,\left\|f_{n} \otimes_{r} f_{n}\right\|_{\mathfrak{S}^{\otimes 2(q-r)}} \rightarrow 0$;
(iv) $\left\|D F_{n}\right\|_{\mathfrak{S}}^{2} \rightarrow q$ in $\mathrm{L}^{2}$.

The implications (i) $\leftrightarrow$ (ii) $\leftrightarrow$ (iii) have been first proved in [36] by means of stochastic calculus techniques. The fact that (iv) is equivalent to either one of conditions (i)-(iii) is proved in [35]. Note that Theorem 3.1 and Proposition 3.2 above provide an alternate proof of the implications (iii) $\rightarrow$ (iv) $\rightarrow$ (i). The implication (ii) $\rightarrow$ (i) can be seen as a drastic simplification of the "method of moments and cumulants", that is a customary tool in order to prove limit theorems for functionals of Gaussian fields (see e.g. [4, 6, 17, 25, 48]). In [38] one can find a multidimensional version of Theorem 3.3.

Remark 3.4 Theorem 3.3 and its generalizations have been applied to a variety frameworks, such as: $p$-variations of stochastic integrals with respect to Gaussian processes [2, 11], quadratic functionals of bivariate Gaussian processes [13], self-intersection local times of fractional Brownian motion [22], approximation schemes for scalar fractional differential equations [30], high-frequency CLTs for random fields on homogeneous spaces [27, 28, 37], needlets analysis on the sphere [1], estimation of self-similarity orders [54], power variations of iterated Brownian motions [33]. We expect that the new bounds proved in Theorem 3.1 and Proposition 3.2 will lead to further refinements of these results. See Section 4 for applications and examples.

As shown in the following statement, the combination of Proposition 3.2 and Theorem 3.3 implies that, on any fixed Wiener chaos, the Kolmogorov, total variation and Wasserstein distances metrize the convergence in distribution towards Gaussian random variables. Other topological characterizations of the set of laws of random variables belonging to a fixed sum of Wiener chaoses are discussed in [23, Ch. VI].

Corollary 3.5 Let the assumptions and notation of Theorem 3.3 prevail. Then, the fact that $F_{n}$ converges in distribution to $Z \sim \mathscr{N}(0,1)$ is equivalent to either one of the following conditions:
(a) $d_{\mathrm{Kol}}\left(F_{n}, Z\right) \rightarrow 0$;
(b) $d_{\mathrm{TV}}\left(F_{n}, Z\right) \rightarrow 0$;
(c) $d_{\mathrm{W}}\left(F_{n}, Z\right) \rightarrow 0$.

Note that the previous result is not trivial, since the topologies induced by $d_{\text {Kol }}, d_{\text {TV }}$ and $d_{\mathrm{W}}$ are stronger than convergence in distribution.

Remark 3.6 (Mehler's formula and Stein's method, I). In [8, Lemma 5.3], Chatterjee has proved the following result (we use a notation which is slightly different from the original statement). Let $Y=g(V)$, where $V=\left(V_{1}, \ldots, V_{n}\right)$ is a vector of centered i.i.d. standard Gaussian random variables, and $g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a smooth function such that: (i) $g$ and its derivatives have subexponential growth at infinity, (ii) $E(g(V))=0$, and (iii) $E\left(g(V)^{2}\right)=1$. Then, for any Lipschitz function $f$, one has that

$$
\begin{equation*}
E[Y f(Y)]=E\left[S(V) f^{\prime}(Y)\right] \tag{3.43}
\end{equation*}
$$

where, for every $v=\left(v_{1}, \ldots, v_{n}\right) \in \mathbb{R}^{n}$,

$$
\begin{equation*}
S(v)=\int_{0}^{1} \frac{1}{2 \sqrt{t}} E\left[\sum_{i=1}^{n} \frac{\partial g}{\partial v_{i}}(v) \frac{\partial g}{\partial v_{i}}(\sqrt{t} v+\sqrt{1-t} V)\right] d t \tag{3.44}
\end{equation*}
$$

so that, for instance, for $Z \sim \mathscr{N}(0,1)$ and by using (1.9), Lemma 1.2 (iii), (1.7) and CauchySchwarz inequality,

$$
\begin{equation*}
d_{\mathrm{TV}}(Y, Z) \leqslant 2 E\left[(S(V)-1)^{2}\right]^{1 / 2} \tag{3.45}
\end{equation*}
$$

We shall prove that (3.43) is a very special case of (1.23). Observe first that, without loss of generality, we can assume that $V_{i}=X\left(h_{i}\right)$, where $X$ is an isonormal process over some Hilbert space of the type $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$ and $\left\{h_{1}, \ldots, h_{n}\right\}$ is an orthonormal system in $\mathfrak{H}$. Since $Y=g\left(V_{1}, \ldots, V_{n}\right)$, we have $D_{a} Y=\sum_{i=1}^{n} \frac{\partial g}{\partial x_{i}}(V) h_{i}(a)$. On the other hand, since $Y$ is centered and square integrable, it admits a chaotic representation of the form $Y=$ $\sum_{q \geqslant 1} I_{q}\left(\psi_{q}\right)$. This implies in particular that $D_{a} Y=\sum_{q=1}^{\infty} q I_{q-1}\left(\psi_{q}(a, \cdot)\right)$. Moreover, one has that $-L^{-1} Y=\sum_{q \geqslant 1} \frac{1}{q} I_{q}\left(\psi_{q}\right)$, so that $-D_{a} L^{-1} Y=\sum_{q \geqslant 1} I_{q-1}\left(\psi_{q}(a, \cdot)\right)$. Now, let $T_{z}, z \geqslant 0$, denote the (infinite dimensional) Ornstein-Uhlenbeck semigroup, whose action on random variables $F \in \mathrm{~L}^{2}(X)$ is given by $T_{z}(F)=\sum_{q \geqslant 0} \mathrm{e}^{-q z} J_{q}(F)$. We can write

$$
\begin{align*}
\int_{0}^{1} \frac{1}{2 \sqrt{t}} T_{\ln (1 / \sqrt{t})}\left(D_{a} Y\right) d t & =\int_{0}^{\infty} e^{-z} T_{z}\left(D_{a} Y\right) d z=\sum_{q \geqslant 1} \frac{1}{q} J_{q-1}\left(D_{a} Y\right) \\
& =\sum_{q \geqslant 1} I_{q-1}\left(\psi_{q}(a, \cdot)\right)=-D_{a} L^{-1} Y \tag{3.46}
\end{align*}
$$

Now recall that Mehler's formula (see e.g. [34, formula (1.54)]) implies that, for every function $f$ with subexponential growth,

$$
T_{z}(f(V))=\left.E\left[f\left(e^{-z} v+\sqrt{1-e^{-2 z}} V\right)\right]\right|_{v=V}, \quad z \geqslant 0
$$

In particular, by applying this last relation to the partial derivatives $\frac{\partial g}{\partial v_{i}}, i=1, \ldots, n$, we deduce from (3.46) that

$$
\int_{0}^{1} \frac{1}{2 \sqrt{t}} T_{\ln (1 / \sqrt{ } t)}\left(D_{a} Y\right) d t=\left.\sum_{i=1}^{n} h_{i}(a) \int_{0}^{1} \frac{1}{2 \sqrt{t}} E\left[\frac{\partial g}{\partial v_{i}}(\sqrt{t} v+\sqrt{1-t} V)\right] d t\right|_{v=V^{*}}
$$

Consequently, (3.43) follows, since

$$
\begin{aligned}
\left\langle D Y,-D L^{-1} Y\right\rangle_{\mathfrak{H}} & =\left\langle\sum_{i=1}^{n} \frac{\partial g}{\partial v_{i}}(V) h_{i},\left.\sum_{i=1}^{n} \int_{0}^{1} \frac{1}{2 \sqrt{t}} E\left[\frac{\partial g}{\partial v_{i}}(\sqrt{t} v+\sqrt{1-t} V)\right] d t\right|_{v=V} h_{i}\right\rangle_{\mathfrak{H}} \\
& =S(V)
\end{aligned}
$$

The following result concerns finite sums of multiple integrals.
Proposition 3.7 For $s \geqslant 2$, fix $s$ integers $2 \leqslant q_{1}<\ldots<q_{s}$. Consider a sequence of the form

$$
Z_{n}=\sum_{i=1}^{s} I_{q_{i}}\left(f_{n}^{i}\right), \quad n \geqslant 1,
$$

where $f_{n}^{i} \in \mathfrak{H}^{\odot q_{i}}$. Set

$$
\mathscr{I}=\left\{(i, j, r) \in\{1, \ldots, s\}^{2} \times \mathbb{N}: 1 \leqslant r \leqslant q_{i} \wedge q_{j} \text { and }\left(r, q_{i}, q_{j}\right) \neq\left(q_{i}, q_{i}, q_{i}\right)\right\} .
$$

Then,

$$
\begin{aligned}
& E\left[\left(1-\left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}}\right)^{2}\right] \leqslant \\
& +2\left(1-\sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2} \\
& +2 s^{2} \sum_{(i, j, r) \in \mathscr{\mathscr { I }}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)! \\
& \\
& \times\left\|f_{n}^{i} \otimes_{q_{i}-r} f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes} \otimes r r}\left\|f_{n}^{j} \otimes_{q_{j}-r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes 2 r}} .
\end{aligned}
$$

In particular, if (as $n \rightarrow \infty) E\left[Z_{n}^{2}\right]=\sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{S}^{\otimes q_{i}}}^{2} \longrightarrow 1$ and if, for any $i=1, \ldots, s$ and $r=1, \ldots, q_{i}-1$, one has that $\left\|f_{n}^{i} \otimes_{r} f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes 2\left(q_{i}-r\right)}} \longrightarrow 0$, then $Z_{n} \xrightarrow{\text { Law }} Z \sim \mathscr{N}(0,1)$ as $n \rightarrow \infty$, and the inequalities in Theorem 3.1 allow to associate bounds with this convergence.

Remark 3.8 1. By using Proposition 3.7, it is possible to prove bounds for limit theorems involving the Gaussian approximation of infinite sums of multiple integrals, such as for instance the CLT proved in [22, Th. 4].
2. Note that, to obtain the convergence result stated in Proposition 3.7, one does not need to suppose that the quantity $E\left[I_{q}\left(f_{i}\right)^{2}\right]=q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H} \otimes q_{i}}^{2}$ is convergent for every $i$. One should compare this finding with the CLTs proved in [38], as well as the Gaussian approximations established in [37].

Proof of Proposition 3.7. Observe first that, without loss of generality, we can assume that $X$ is an isonormal process over some Hilbert space of the type $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$. For every $a \in A$, it is immediately checked that

$$
D_{a} Z_{n}=\sum_{i=1}^{s} q_{i} I_{q_{i}-1}\left(f_{n}^{i}(\cdot, a)\right)
$$

and

$$
-D_{a}\left(L^{-1} Z_{n}\right)=D_{a}\left(\sum_{i=1}^{s} \frac{1}{q_{i}} I_{q_{i}}\left(f_{n}^{i}\right)\right)=\sum_{i=1}^{s} I_{q_{i}-1}\left(f_{n}^{i}(\cdot, a)\right)
$$

This yields, using in particular the multiplication formula (2.29):

$$
\begin{aligned}
& \left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}} \\
= & \sum_{i, j=1}^{s} q_{i} \int_{A} I_{q_{i}-1}\left(f_{n}^{i}(\cdot, a)\right) I_{q_{j}-1}\left(f_{n}^{j}(\cdot, a)\right) \mu(d a) \\
= & \sum_{i, j=1}^{s} q_{i} \sum_{r=0}^{q_{i} \wedge q_{j}-1} r!\binom{q_{i}-1}{r}\binom{q_{j}-1}{r} I_{q_{i}+q_{j}-2-2 r}\left(\int_{A} f_{n}^{i}(\cdot, a) \otimes_{r} f_{n}^{j}(\cdot, a) \mu(d a)\right) \\
= & \sum_{i, j=1}^{s} q_{i} \sum_{r=0}^{q_{i} \wedge q_{j}-1} r!\binom{q_{i}-1}{r}\binom{q_{j}-1}{r} I_{q_{i}+q_{j}-2-2 r}\left(f_{n}^{i} \otimes_{r+1} f_{n}^{j}\right) \\
= & \sum_{i, j=1}^{s} q_{i} \sum_{r=1}^{q_{i} \wedge q_{j}}(r-1)!\binom{q_{i}-1}{r-1}\binom{q_{j}-1}{r-1} I_{q_{i}+q_{j}-2 r}\left(f_{n}^{i} \otimes_{r} f_{n}^{j}\right) \\
= & \sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}}^{2} \otimes q_{i}
\end{aligned} \sum_{(i, j, r) \in \mathscr{I}} q_{i}(r-1)!\binom{q_{i}-1}{r-1}\binom{q_{j}-1}{r-1} I_{q_{i}+q_{j}-2 r}\left(f_{n}^{i} \otimes_{r} f_{n}^{j}\right) . .
$$

Thus, by using (among others) inequalities of the type $\left(a_{1}+\ldots+a_{v}\right)^{2} \leqslant v\left(a_{1}^{2}+\ldots+a_{v}^{2}\right)$, the isometric properties of multiple integrals as well $\left\|f \widetilde{\otimes}_{r} g\right\| \leqslant\left\|f \otimes_{r} g\right\|$, we obtain

$$
\begin{aligned}
& E\left(\left[\left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}}-1\right]^{2}\right) \\
& \leqslant 2\left(1-\sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2} \\
&+2 E\left(\sum_{(i, j, r) \in \mathscr{\mathscr { I }}} q_{i}(r-1)!\binom{q_{i}-1}{r-1}\binom{q_{j}-1}{r-1} I_{q_{i}+q_{j}-2 r}\left(f_{n}^{i} \otimes_{r} f_{n}^{j}\right)\right)^{2} \\
& \leqslant 2\left(1-\sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2} \\
&+2 s^{2} \sum_{(i, j, r) \in \mathscr{\mathscr { I }}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)!\left\|f_{n}^{i} \otimes_{r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes q_{i}+\mathscr{G}_{j}-2 r}}^{2} \\
& \leqslant 2\left(1-\sum_{i=1}^{s} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{S}^{\otimes q_{i}}}^{2}\right)^{2} \\
&+2 s^{2} \sum_{(i, j, r) \in \mathscr{I}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)! \\
& \times\left\|f_{n}^{i} \otimes_{q_{i}-r} f_{n}^{i}\right\|_{\mathfrak{S} \otimes 2 r}\left\|f_{n}^{j} \otimes_{q_{j}-r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes 2 r}},
\end{aligned}
$$

the last inequality being a consequence of the (easily verified) relation

$$
\left\|f_{n}^{i} \otimes_{r} f_{n}^{j}\right\|_{\mathfrak{S}^{\otimes q_{i}+q_{j}-2 r}}^{2}=\left\langle f_{n}^{i} \otimes_{q_{i}-r} f_{n}^{i}, f_{n}^{j} \otimes_{q_{j}-r} f_{n}^{j}\right\rangle_{\mathfrak{S}^{\otimes 2 r}}
$$

### 3.2 A property of $\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}$

Before dealing with Gamma approximations, we shall prove the a.s. positivity of a specific projection of the random variable $\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}}$ appearing in Theorem 3.1. This fact will be used in the proof of the main result of the next section.

Proposition 3.9 Let $F \in \mathbb{D}^{1,2}$. Then, $P$-a.s.,

$$
\begin{equation*}
E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} \mid F\right] \geqslant 0 . \tag{3.47}
\end{equation*}
$$

Proof. Let $g$ be a non-negative real function, and set $G(x)=\int_{0}^{x} g(t) d t$, with the usual convention $\int_{0}^{x}=-\int_{x}^{0}$ for $x<0$. Since $G$ is increasing and vanishing at zero, we have $x G(x) \geqslant 0$ for all $x \in \mathbb{R}$. In particular, $E(F G(F)) \geqslant 0$. Moreover,

$$
E[F G(F)]=E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} g(F)\right]=E\left[E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} \mid F\right] g(F)\right] .
$$

We therefore deduce that

$$
E\left[E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} \mid F\right] \mathbf{1}_{A}\right] \geqslant 0
$$

for any $\sigma(F)$-measurable set $A$. This implies the desired conclusion.

Remark 3.10 According to Goldstein and Reinert [19], for $F$ as in the previous statement, there exists a random variable $F^{*}$ having the $F$-zero biased distribution, that is, $F^{*}$ is such that, for every absolutely continuous function $f$,

$$
E\left[f^{\prime}\left(F^{*}\right)\right]=E[F f(F)]
$$

By the computations made in the previous proof, one also has that

$$
E\left[g\left(F^{*}\right)\right]=E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} g(F)\right]
$$

for any real-valued and smooth function $g$. This implies, in particular, that the conditional expectation $E\left[\left\langle D F,-D L^{-1} F\right\rangle_{\mathfrak{H}} \mid F\right]$ is a version of the Radon-Nikodym derivative of the law of $F^{*}$ with respect to the law of $F$, whenever the two laws are equivalent.

### 3.3 Gamma approximations

We now combine Malliavin calculus with the Gamma approximations discussed in the second part of Section 1.2.

Theorem 3.11 Fix $\nu>0$ and let $F(\nu)$ have a centered Gamma distribution with parameter $\nu$. Let $G \in \mathbb{D}^{1,2}$ be such that $E(G)=0$ and the law of $G$ is absolutely continuous with respect to the Lebesgue measure. Then:

$$
\begin{align*}
d_{\mathscr{H}_{1}}(G, F(\nu)) & \leqslant K_{1} E\left[\left(2 \nu+2 G-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2}  \tag{3.48}\\
d_{\mathscr{H}_{2}}(G, F(\nu)) & \leqslant K_{2} E\left[\left(2 \nu+2 G-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2} \tag{3.49}
\end{align*}
$$

where $\mathscr{H}_{1}$ and $\mathscr{H}_{2}$ are defined in (1.14)-(1.15), $K_{1}:=\max \left\{1,1 / \nu+2 / \nu^{2}\right\}$ and $K_{2}:=\max \{\sqrt{2 \pi / \nu}$, $\left.1 / \nu+2 / \nu^{2}\right\}$.

Proof. We will only prove (3.48), the proof of (3.49) being analogous. Fix $\nu>0$. Thanks to (1.20) and (1.23) (in the case $Y=G$ ) and by applying Cauchy-Schwarz, we deduce that

$$
\begin{aligned}
d_{\mathscr{H}_{1}}(G, F(\nu)) & \leqslant \sup _{\mathscr{F}_{1}} \mid E\left[f^{\prime}(G)\left(2(\nu+G)_{+}-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right) \mid\right. \\
& \leqslant K_{1} \times E\left[\left(2(\nu+G)_{+}-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2} \\
& \leqslant K_{1} \times E\left[\left(2(\nu+G)-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right)^{2}\right]^{1 / 2}
\end{aligned}
$$

where the last inequality is a consequence of the fact that $E\left[\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}} \mid G\right] \geqslant 0$ (thanks to Proposition 3.9).

Remark 3.12 (Mehler's formula and Stein's method, II). Define $Y=g(V)$ as in Remark 3.6. Then, since (3.44) and (3.45) are in order, one deduces from Theorem 3.11 that, for every $\nu>0$,

$$
d_{\mathscr{H}_{1}}(Y, F(\nu)) \leqslant K_{1} E\left[(2 \nu+2 Y-S(V))^{2}\right]^{1 / 2}
$$

An analogous estimate holds for $d_{\mathscr{H}_{2}}$, when applied to the case where $\nu \geqslant 1$ is an integer.

We will now connect the previous results to the main findings of [32]. To do this, we shall provide explicit estimates of the bounds appearing in Theorem 3.11, in the case where $G$ belongs to a fixed Wiener chaos of even order $q$.

Proposition 3.13 Let $q \geqslant 2$ be an even integer, and let $G=I_{q}(g)$, where $g \in \mathfrak{H}^{\odot q}$. Then,

$$
\begin{align*}
& E\left[\left(2 \nu+2 G-\left\langle D G,-D L^{-1} G\right\rangle_{\mathfrak{H}}\right)^{2}\right]=E\left[\left(2 \nu+2 G-q^{-1}\|D G\|_{\mathfrak{H}}^{2}\right)^{2}\right]  \tag{3.50}\\
& \leqslant\left(2 \nu-q!\|g\|_{\mathfrak{H} \otimes q}^{2}\right)^{2}+  \tag{3.51}\\
& \quad+q^{2} \sum_{\substack{r \in\{1, \ldots, q-1\} \\
r \neq q / 2}}(2 q-2 r)!(r-1)!^{2}\binom{q-1}{r-1}^{4}\left\|g \otimes_{r} g\right\|_{\mathfrak{H}^{\otimes 2(q-r)}}^{2}+ \\
&  \tag{3.52}\\
& \quad+4 q!\left\|c_{q}^{-1} \times g \widetilde{\otimes}_{q / 2} g-g\right\|_{\mathfrak{H}^{\otimes q}}^{2}
\end{align*}
$$

where

$$
\begin{equation*}
c_{q}:=\frac{1}{(q / 2)!\binom{q-1}{q / 2-1}^{2}}=\frac{4}{(q / 2)!\binom{q}{q / 2}^{2}} . \tag{3.53}
\end{equation*}
$$

Proof. By using (3.42) we deduce that

$$
\begin{aligned}
q^{-1}\|D G\|_{\mathfrak{H}}^{2}-2 \nu-2 G & =\left(q!\|g\|_{\mathfrak{H}^{\otimes q}}^{2}-2 \nu\right)+ \\
& +q \sum_{\substack{r \in\{1, \ldots, q-1\} \\
r \neq q / 2}}(r-1)!\binom{q-1}{r-1}^{2} I_{2(q-r)}\left(g \widetilde{\otimes}_{r} g\right)+ \\
& +q(q / 2-1)!\binom{q-1}{q / 2-1} I_{q}\left(g \widetilde{\otimes}_{q / 2} g-2 g\right) .
\end{aligned}
$$

The conclusion is obtained by using the isometric properties of multiple Wiener-Itô integrals, as well as the relation $\left\|g \widetilde{\otimes}_{r} g\right\|_{\mathfrak{H}^{\otimes 2(q-r)}} \leqslant\left\|g \otimes_{r} g\right\|_{\mathfrak{H}^{\otimes 2(q-r)}}$, for every $r \in\{1, \ldots, q-1\}$ such that $r \neq q / 2$.

By using Proposition 3.13, we immediately recover the implications (iv) $\rightarrow$ (iii) $\rightarrow$ (i) in the statement of the following result, recently proved in [32, Th. 1.2].

Theorem 3.14 ([32]) Let $\nu>0$ and let $F(\nu)$ have a centered Gamma distribution with parameter $\nu$. Fix an even integer $q \geqslant 2$, and define $c_{q}$ according to (3.53). Consider a sequence of the type $G_{n}=I_{q}\left(g_{n}\right)$, where $n \geqslant 1$ and $g_{n} \in \mathfrak{H}^{\odot q}$, and suppose that

$$
\lim _{n \rightarrow \infty} E\left[G_{n}^{2}\right]=\lim _{n \rightarrow \infty} q!\left\|g_{n}\right\|_{\mathfrak{H}^{\otimes q}}^{2}=2 \nu
$$

Then, the following four conditions are equivalent:
(i) as $n \rightarrow \infty$, the sequence $\left(G_{n}\right)_{n \geqslant 1}$ converges in distribution to $F(\nu)$;
(ii) $\lim _{n \rightarrow \infty} E\left[G_{n}^{4}\right]-12 E\left[G_{n}^{3}\right]=12 \nu^{2}-48 \nu$;
(iii) as $n \rightarrow \infty,\left\|D G_{n}\right\|_{\mathfrak{H}}^{2}-2 q G_{n} \longrightarrow 2 q \nu$ in $\mathrm{L}^{2}$.
(iv) $\lim _{n \rightarrow \infty}\left\|g_{n} \widetilde{\otimes}_{q / 2} g_{n}-c_{q} \times g_{n}\right\|_{\mathfrak{H}^{\otimes q}}=0$, where $c_{q}$ is given by (3.53), and $\lim _{n \rightarrow \infty} \| g_{n} \otimes_{r}$ $g_{n} \|_{\mathfrak{H}^{\otimes 2(q-r)}}=0$, for every $r=1, \ldots, q-1$ such that $r \neq q / 2$.

Observe that $E\left(F(\nu)^{2}\right)=2 \nu, E\left(F(\nu)^{3}\right)=8 \nu$ and $E\left(F(\nu)^{4}\right)=48 \nu+12 \nu^{2}$, so that the implication (ii) $\rightarrow$ (i) in the previous statement can be seen as a further simplification of the method of moments and cumulants, as applied to non-central limit theorems (see e.g. [49], and the references therein, for a survey of classic non-central limit theorems). Also, the combination of Proposition 3.13 and Theorem 3.14 shows that, inside a fixed Wiener chaos of even order, one has that: (i) $d_{\mathscr{H}_{1}}$ metrizes the weak convergence towards centered Gamma distributions, and (ii) $d_{\mathscr{H}_{2}}$ metrizes the weak convergence towards centered $\chi_{2}$ distributions with arbitrary degrees of freedom.

The following result concerns the Gamma approximation of a sum of two multiple integrals. Note, at the cost of a quite heavy notation, one could easily establish analogous estimates for sums of three or more integrals. The reader should compare this result with Proposition 3.7.

Proposition 3.15 Fix two real numbers $\nu_{1}, \nu_{2}>0$, as well as two even integers $2 \leqslant q_{1}<q_{2}$. Set $\nu=\nu_{1}+\nu_{2}$ and suppose (for the sake of simplicity) that $q_{2}>2 q_{1}$. Consider a sequence of the form

$$
Z_{n}=I_{q_{1}}\left(f_{n}^{1}\right)+I_{q_{2}}\left(f_{n}^{2}\right), \quad n \geqslant 1
$$

where $f_{n}^{i} \in \mathfrak{H}^{\odot q_{i}}$. Set

$$
\mathscr{J}=\left\{(i, j, r) \in\{1,2\}^{2} \times \mathbb{N}: 1 \leqslant r \leqslant q_{i} \wedge q_{j} \text { and, whenever } i=j, r \neq q_{i} \text { and } r \neq \frac{q_{i}}{2}\right\} .
$$

Then

$$
\begin{align*}
& E\left[\left(2 Z_{n}+2 \nu-\left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}}\right)^{2}\right] \\
& \leqslant 3\left(2 \nu-\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2}+24 \sum_{i=1,2} c_{q_{i}}^{-2} q_{i}!\left\|f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2} \\
& +12 \sum_{(i, j, r) \in \mathscr{J}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)!  \tag{3.54}\\
& \times\left\|f_{n}^{i} \otimes_{q_{i}-r} f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes 2 r}}\left\|f_{n}^{j} \otimes_{q_{j}-r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes 2 r}} .
\end{align*}
$$

In particular, if
(i) $E\left[Z_{n}^{2}\right]=\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2} \longrightarrow 2 \nu$ as $n \rightarrow \infty$,
(ii) for $i=1,2,\left\|f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}} \longrightarrow 0$ as $n \rightarrow \infty$, where $c_{q_{i}}$ is defined in Theorem 3.14,
(iii) for any $i=1,2$ and $r=1, \ldots, q_{i}-1$ such that $r \neq \frac{q_{i}}{2},\left\|f_{n}^{i} \otimes_{r} f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes 2\left(q_{i}-r\right)}} \longrightarrow 0$ as $n \rightarrow \infty$,
then $Z_{n} \xrightarrow{\text { Law }} F(\nu)$ as $n \rightarrow \infty$, and the combination of Theorem 3.1 and (3.54) allows to associate explicit bounds with this convergence.

Proof of Proposition 3.15. We have (see the proof of Proposition 3.7)

$$
\begin{aligned}
& \left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}}-2 Z_{n}-2 \nu \\
& =\left(\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}-2 \nu\right)+\sum_{i=1,2} 2 c_{q_{i}}^{-1} I_{q_{i}}\left(f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right) \\
& \quad+\sum_{(i, j, r) \in \mathscr{J}} q_{i}(r-1)!\binom{q_{i}-1}{r-1}\binom{q_{j}-1}{r-1} I_{q_{i}+q_{j}-2 r}\left(f_{n}^{i} \otimes_{r} f_{n}^{j}\right) .
\end{aligned}
$$

Thus

$$
\begin{aligned}
& E\left(\left[\left\langle D Z_{n},-D L^{-1} Z_{n}\right\rangle_{\mathfrak{H}}-2 Z_{n}-2 \nu\right]^{2}\right) \\
& \leqslant 3\left(2 \nu-\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2}+24 \sum_{i=1,2} c_{q_{i}}^{-2} q_{i}!\left\|f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2} \\
& +3 E\left(\sum_{(i, j, r) \in \mathscr{J}} q_{i}(r-1)!\binom{q_{i}-1}{r-1}\binom{q_{j}-1}{r-1} I_{q_{i}+q_{j}-2 r}\left(f_{n}^{i} \otimes_{r} f_{n}^{j}\right)\right)^{2} \\
& \leqslant 3\left(2 \nu-\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2}+24 \sum_{i=1,2} c_{q_{i}}^{-2} q_{i}!\left\|f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2} \\
& +12 \sum_{(i, j, r) \in \mathscr{J}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)!\left\|f_{n}^{i} \otimes_{r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes q_{i}+q_{j}-2 r}}^{2} \\
& \leqslant 3\left(2 \nu-\sum_{i=1,2} q_{i}!\left\|f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2}\right)^{2}+24 \sum_{i=1,2} c_{q_{i}}^{-2} q_{i}!\left\|f_{n}^{i} \widetilde{\otimes}_{q_{i} / 2} f_{n}^{i}-c_{q_{i}} \times f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes q_{i}}}^{2} \\
& +12 \sum_{(i, j, r) \in \mathscr{F}} q_{i}^{2}(r-1)!^{2}\binom{q_{i}-1}{r-1}^{2}\binom{q_{j}-1}{r-1}^{2}\left(q_{i}+q_{j}-2 r\right)! \\
& \times\left\|f_{n}^{i} \otimes_{q_{i}-r} f_{n}^{i}\right\|_{\mathfrak{H}^{\otimes 2 r}}\left\|f_{n}^{j} \otimes_{q_{j}-r} f_{n}^{j}\right\|_{\mathfrak{H}^{\otimes 2 r}} .
\end{aligned}
$$

## 4 Berry-Esséen bounds in the Breuer-Major CLT

In this section, we use our main results in order to derive an explicit Berry-Esséen bound for the celebrated Breuer-Major CLT for Gaussian-subordinated random sequences. For simplicity, we focus on sequences that can be represented as Hermite-type functions of the (normalized) increments of a fractional Brownian motion. Our framework include examples of Gaussian sequences whose autocovariance functions display long dependence. Plainly, the techniques developed in this paper can also accommodate the analysis of more general transformations (for instance, obtained from functions with an arbitrary Hermite rank - see [50]), as well as alternative covariance structures.

### 4.1 General setup

We recall that a fractional Brownian motion (fBm) $B=\left\{B_{t}: t \in[0,1]\right\}$, of Hurst index $H \in(0,1)$, is a centered Gaussian process, started from zero and with covariance function $E\left(B_{s} B_{t}\right)=R_{H}(s, t)$, where

$$
R_{H}(s, t)=\frac{1}{2}\left(t^{2 H}+s^{2 H}-|t-s|^{2 H}\right) ; \quad s, t \in[0,1]
$$

If $H=1 / 2$, then $R_{H}(s, t)=\min (s, t)$ and $B$ is a standard Brownian motion. For any choice of the Hurst parameter $H \in(0,1)$, the Gaussian space generated by $B$ can be identified with
an isonormal Gaussian process of the type $X=\{X(h): h \in \mathfrak{H}\}$, where the real and separable Hilbert space $\mathfrak{H}$ is defined as follows: (i) denote by $\mathscr{E}$ the set of all $\mathbb{R}$-valued step functions on $[0,1]$, (ii) define $\mathfrak{H}$ as the Hilbert space obtained by closing $\mathscr{E}$ with respect to the scalar product

$$
\left\langle\mathbf{1}_{[0, t]}, \mathbf{1}_{[0, s]}\right\rangle_{\mathfrak{H}}=R_{H}(t, s) .
$$

In particular, with such a notation one has that $B_{t}=X\left(\mathbf{1}_{[0, t]}\right)$. Note that, if $H=1 / 2$, then $\mathfrak{H}=L^{2}[0,1]$; when $H>1 / 2$, the space $\mathfrak{H}$ coincides with the space of distributions $f$ such that $s^{\frac{1}{2}-H} \mathscr{I}_{0+}^{H-\frac{1}{2}}\left(f(u) u^{H-\frac{1}{2}}\right)(s)$ belongs to $L^{2}[0,1]$; when $H<1 / 2$ one has that $\mathfrak{H}$ is $\mathscr{I}_{0+}^{H-\frac{1}{2}}\left(L^{2}[0,1]\right)$. Here, $\mathscr{I}_{0+}^{H-\frac{1}{2}}$ denotes the action of the fractional Riemann-Liouville operator, defined as

$$
\mathscr{I}_{0+}^{H-\frac{1}{2}} f(x)=\frac{1}{\Gamma\left(H-\frac{1}{2}\right)} \int_{0}^{x}(x-y)^{H-\frac{3}{2}} f(y) d y .
$$

The reader is referred e.g. to [34] for more details on fBm and fractional operators.

### 4.2 A Berry-Esséen bound

In what follows, we will be interested in the asymptotic behaviour (as $n \rightarrow \infty$ ) of random vectors that are subordinated to the array

$$
\begin{equation*}
V_{n, H}=\left\{n^{H}\left(B_{(k+1) / n}-B_{k / n}\right): k=0, \ldots, n-1\right\}, n \geqslant 1 \tag{4.55}
\end{equation*}
$$

Note that, for every $n \geqslant 1$, the law of $V_{n, H}$ in (4.55) coincides with the law of the first $n$ instants of a centered stationary Gaussian sequence indexed by $\{0,1,2, \ldots\}$ and with autocovariance function given by

$$
\rho_{H}(k)=\frac{1}{2}\left(|k+1|^{2 H}+|k-1|^{2 H}-2|k|^{2 H}\right), \quad k \in \mathbb{Z}
$$

(in particular, $\rho_{H}(0)=1$ and $\rho_{H}(k)=\rho_{H}(-k)$ ). From this last expression, one deduces that the components of the vector $V_{n, H}$ are: (a) i.i.d. for $H=1 / 2$, (b) negatively correlated for $H \in(0,1 / 2)$ and (c) positively correlated for $H \in(1 / 2,1)$. In particular, if $H \in(1 / 2,1)$, then $\sum_{k} \rho_{H}(k)=+\infty$ : in this case, one customarily says that $\rho_{H}$ exhibits long-range dependence (or, equivalently, long memory - see e.g. [52] for a general discussion of this point).
Now denote by $H_{q}, q \geqslant 2$, the $q$ th Hermite polynomial, defined as

$$
H_{q}(x)=\frac{(-1)^{q}}{q!} e^{\frac{x^{2}}{2}} \frac{d^{q}}{d x^{q}} e^{-\frac{x^{2}}{2}}, x \in \mathbb{R}
$$

For instance, $H_{2}(x)=\left(x^{2}-1\right) / 2, H_{3}(x)=\left(x^{3}-3 x\right) / 6$, and so on. Finally, set

$$
\sigma=\sqrt{\frac{1}{q!} \sum_{t \in \mathbb{Z}} \rho_{H}(t)^{q}}
$$

and define

$$
\begin{equation*}
Z_{n}=\frac{1}{\sigma \sqrt{n}} \sum_{k=0}^{n-1} H_{q}\left(n^{H}\left(B_{(k+1) / n}-B_{k / n}\right)\right)=\frac{n^{q H-\frac{1}{2}}}{q!\sigma} \sum_{k=0}^{n-1} I_{q}\left(\delta_{k / n}^{\otimes q}\right) \tag{4.56}
\end{equation*}
$$

where $I_{q}$ denotes the $q$ th multiple integral with respect to the isonormal process associated with $B$ (see Section 2). For simplicity, here (and for the rest of this section) we write $\delta_{k / n}$ instead of $\mathbf{1}_{[k / n,(k+1) / n]}$, and also $\delta_{k / n}^{\otimes q}=\delta_{k / n} \otimes \cdots \otimes \delta_{k / n}$ ( $q$ times). Note that in (4.56) we have used the standard relation: $q!H_{q}(h)=I_{q}\left(h^{\otimes q}\right)$ for every $h \in \mathfrak{H}$ such that $\|h\|_{\mathfrak{H}}=1$ (see e.g. [34, Ch. 1]).

Now observe that, for every $q \geqslant 2$, one has that $\sum_{t}\left|\rho_{H}(t)\right|^{q}<\infty$ if, and only if, $H \in$ $\left(0, \frac{2 q-1}{2 q}\right)$. Moreover, in this case, $E\left(Z_{n}^{2}\right) \rightarrow 1$ as $n \rightarrow \infty$. As a consequence, according to Breuer and Major's well-known result [4, Theorem 1], as $n \rightarrow \infty$

$$
Z_{n} \rightarrow Z \sim \mathscr{N}(0,1) \quad \text { in distribution. }
$$

To the authors' knowledge, the following statement contains the first Berry-Esséen bound ever proved for the Breuer-Major CLT:

Theorem 4.1 As $n \rightarrow \infty, Z_{n}$ converges in law towards $Z \sim \mathscr{N}(0,1)$. Moreover, there exists a constant $c_{H}$, depending uniquely on $H$, such that, for any $n \geqslant 1$ :

$$
\sup _{z \in \mathbb{R}}\left|P\left(Z_{n} \leqslant z\right)-P(Z \leqslant z)\right| \leqslant c_{H} \times \begin{cases}n^{-\frac{1}{2}} & \text { if } H \in\left(0, \frac{1}{2}\right] \\ n^{H-1} & \text { if } H \in\left(\frac{1}{2}, \frac{2 q-3}{2 q-2}\right] \\ n^{q H-q+\frac{1}{2}} & \text { if } H \in\left(\frac{2 q-3}{2 q-2}, \frac{2 q-1}{2 q}\right)\end{cases}
$$

Remark 4.2 1. Theorem 1.6 (see the Introduction) can be proved by simply setting $q=2$ in Theorem 4.1. Observe that in this case one has $\frac{2 q-3}{2 q-2}=\frac{1}{2}$, so that the middle line in the previous display becomes immaterial.
2. When $H>\frac{2 q-1}{2 q}$, the sequence $Z_{n}$ does not converge in law towards a Gaussian random variable. Indeed, in this case a non-central limit theorem takes place. See e.g. Nourdin et al. [31, Th. 1] or Taqqu $[50,51]$.
3. As discussed in [4, p. 429], it is in general not possible to derive CLTs such as the one in Theorem 4.1 from mixing-type conditions. In particular, it seems unfeasible to deduce Theorem 4.1 from any mixing characterization of the increments of fractional Brownian motion (as the one proved e.g. by Picard in [39, Theorem A.1]). See e.g. Tikhomirov [53] for general derivations of Berry-Esséen bounds from strong mixing conditions.

### 4.3 Proof of Theorem 4.1

We have

$$
D Z_{n}=\frac{n^{q H-\frac{1}{2}}}{(q-1)!\sigma} \sum_{k=0}^{n-1} I_{q-1}\left(\delta_{k / n}^{\otimes q-1}\right) \delta_{k / n}
$$

hence

$$
\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}=\frac{n^{2 q H-1}}{(q-1)!^{2} \sigma^{2}} \sum_{k, \ell=0}^{n-1} I_{q-1}\left(\delta_{k / n}^{\otimes q-1}\right) I_{q-1}\left(\delta_{\ell / n}^{\otimes q-1}\right)\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}
$$

By the multiplication formula (2.29):

$$
I_{q-1}\left(\delta_{k / n}^{\otimes q-1}\right) I_{q-1}\left(\delta_{\ell / n}^{\otimes q-1}\right)=\sum_{r=0}^{q-1} r!\binom{q-1}{r}^{2} I_{2 q-2-2 r}\left(\delta_{k / n}^{\otimes q-1-r} \widetilde{\otimes} \delta_{\ell / n}^{q-1-r}\right)\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{r}
$$

Consequently,

$$
\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}=\frac{n^{2 q H-1}}{(q-1)!^{2} \sigma^{2}} \sum_{r=0}^{q-1} r!\binom{q-1}{r}^{2} \sum_{k, \ell=0}^{n-1} I_{2 q-2-2 r}\left(\delta_{k / n}^{\otimes q-1-r} \widetilde{\otimes} \delta_{\ell / n}^{q-1-r}\right)\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{r+1}
$$

Thus, we can write

$$
\frac{1}{q}\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}-1=\sum_{r=0}^{q-1} A_{r}(n)-1
$$

where

$$
A_{r}(n)=\frac{r!\binom{q-1}{r}^{2}}{q(q-1)!^{2} \sigma^{2}} n^{2 q H-1} \sum_{k, \ell=0}^{n-1} I_{2 q-2-2 r}\left(\delta_{k / n}^{\otimes q-1-r} \widetilde{\otimes} \delta_{\ell / n}^{q-1-r}\right)\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{r+1}
$$

We will need the following easy Lemma (the proof is omitted). Here and for the rest of the proof of Theorem 4.1, the notation $a_{n} \asymp b_{n}$ means that $\sup _{n \geqslant 1}\left|a_{n}\right| /\left|b_{n}\right|<\infty$.
Lemma 4.3 1. We have $\rho_{H}(n) \asymp|n|^{2 H-2}$.
2. For any $\alpha \in \mathbb{R}$, we have

$$
\sum_{k=1}^{n-1} k^{\alpha} \asymp 1+n^{\alpha+1}
$$

3. If $\alpha \in(-\infty,-1)$, we have

$$
\sum_{k=n}^{\infty} k^{\alpha} \asymp n^{\alpha+1}
$$

By using elementary computations (in particular, observe that $n^{2 H}\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}=\rho_{H}(k-\ell)$ ) and then Lemma 4.3, it is easy to check that

$$
\begin{aligned}
A_{q-1}(n)-1 & =\frac{1}{q!\sigma^{2}} n^{2 q H-1} \sum_{k, \ell=0}^{n-1}\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{q}-1 \\
& =\frac{1}{q!\sigma^{2}}\left(\frac{1}{n} \sum_{k, \ell=0}^{n-1} \rho_{H}(k-\ell)^{q}-\sum_{t \in \mathbb{Z}} \rho_{H}(t)^{q}\right) \\
& =\frac{1}{q!\sigma^{2}}\left(\frac{1}{n} \sum_{|t|<n}(n-1+t) \rho_{H}(t)^{q}-\sum_{t \in \mathbb{Z}} \rho_{H}(t)^{q}\right) \\
& =\frac{1}{q!\sigma^{2}}\left(\frac{1}{n} \sum_{|t|<n}(t-1) \rho_{H}(t)^{q}-\sum_{|t| \geqslant n} \rho_{H}(t)^{q}\right) \\
& \asymp \frac{1}{n} \sum_{t=1}^{n-1} t^{2 q H-2 q+1}+\sum_{t=n}^{\infty} t^{2 q H-2 q} \asymp n^{-1}+n^{2 q H-2 q+1} .
\end{aligned}
$$

Now, we assume that $r \leqslant q-2$ is fixed. We have

$$
\begin{aligned}
E\left|A_{r}(n)\right|^{2}= & c(H, r, q) n^{4 q H-2} \sum_{i, j, k, \ell=0}^{n-1}\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{r+1}\left\langle\delta_{i / n}, \delta_{j / n}\right\rangle_{\mathfrak{H}}^{r+1} \\
& \times\left\langle\delta_{k / n}^{\otimes q-1-r} \widetilde{\otimes} \delta_{\ell / n}^{q-1-r}, \delta_{i / n}^{\otimes q-1-r} \widetilde{\otimes} \delta_{j / n}^{q-1-r}\right\rangle_{\mathfrak{H} \otimes 2 q-2-2 r} \\
= & \sum_{\substack{\alpha, \beta \geqslant 0 \\
\alpha+\beta=q-r-1}} \sum_{\substack{\gamma, \delta, \delta \geqslant 0 \\
\gamma+\delta=q-r-1}} c(H, r, q, \alpha, \beta, \gamma, \delta) B_{r, \alpha, \beta, \gamma, \delta}(n)
\end{aligned}
$$

where $c(\cdot)$ denotes a generic constant depending only on the objects inside its argument, and

$$
\begin{aligned}
& B_{r, \alpha, \beta, \gamma, \delta}(n)=n^{4 q H-2} \sum_{i, j, k, \ell=0}^{n-1}\left\langle\delta_{k / n}, \delta_{\ell / n}\right\rangle_{\mathfrak{H}}^{r+1}\left\langle\delta_{i / n}, \delta_{j / n}\right\rangle_{\mathfrak{H}}^{r+1}\left\langle\delta_{k / n}, \delta_{i / n}\right\rangle_{\mathfrak{H}}^{\alpha} \\
& \times\left\langle\delta_{k / n}, \delta_{j / n}\right\rangle_{\mathfrak{H}}^{\beta}\left\langle\delta_{\ell / n}, \delta_{i / n}\right\rangle_{\mathfrak{H}}^{\gamma}\left\langle\delta_{\ell / n}, \delta_{j / n}\right\rangle_{\mathfrak{H}}^{\delta} \\
&=n^{-2} \sum_{i, j, k, \ell=0}^{n-1} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \\
& \times \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} .
\end{aligned}
$$

When $\alpha, \beta, \gamma, \delta$ are fixed, we can decompose the sum $\sum_{i, j, k, \ell}$ appearing in $B_{r, \alpha, \beta, \gamma, \delta}(n)$ just above, as follows:

$$
\begin{aligned}
& \sum_{i=j=k=\ell}+\left(\sum_{\substack{i=j=k \\
\ell \neq i}}+\sum_{\substack{i=j=\ell \\
k \neq i}}+\sum_{\substack{i=k=\ell \\
j \neq i}}+\sum_{\substack{j=k=\ell \\
i \neq j}}\right)+\left(\sum_{\substack{i=j, k=\ell \\
k \neq i}}+\sum_{\substack{i=k, j=\ell \\
j \neq i}}+\sum_{\substack{i=\ell, j=k \\
j \neq i}}\right) \\
& +\left(\sum_{\substack{i=j, k \neq i \\
k \neq \ell, \ell \neq i}}+\sum_{\substack{i=k, j \neq i \\
j \neq \ell, k \neq \ell}}+\sum_{\substack{i=\ell, k \neq i \\
k \neq j, j \neq i}}+\sum_{\substack{j=k, k \neq i \\
k \neq \ell, \ell \neq i}}+\sum_{\substack{j=\ell, k \neq i \\
k \neq \ell, \ell \neq i}}+\sum_{\substack{k=\ell, k \neq i \\
k \neq j, j \neq i}}\right)+\sum_{\text {are all different }}
\end{aligned}
$$

(all these sums must be understood as being defined over indices $\{i, j, k, \ell\} \in\{0, \ldots, n-1\}^{4}$ ). Now, we will deal with each of these fifteen sums separately.

The first sum is particularly easy to handle: indeed, it is immediately checked that

$$
n^{-2} \sum_{i=j=k=\ell} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \asymp n^{-1} .
$$

For the second sum, can write

$$
\begin{aligned}
& n^{-2} \sum_{\substack{i=j=k \\
\ell \neq i}} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \\
\asymp & n^{-2} \sum_{i \neq \ell} \rho_{H}(\ell-i)^{q} \asymp n^{-1} \sum_{\ell=1}^{n-1} \ell^{2 q H-2 q}=n^{-1}+n^{2 q H-2 q} \quad \text { by Lemma 4.3. }
\end{aligned}
$$

For the third sum, we can proceed analogously and we also obtain $n^{-1}+n^{2 q H-2 q}$ for bound.

For the fourth sum, we write

$$
\begin{aligned}
& n^{-2} \sum_{\substack{i=k=\ell \\
j \neq i}} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \\
\asymp & n^{-2} \sum_{i \neq j} \rho_{H}(j-i)^{r+1+\beta+\delta} \asymp n^{-2} \sum_{i \neq j}|j-i|^{(r+1+\beta+\delta)(2 H-2)} \asymp n^{-2} \sum_{i \neq j}|j-i|^{2 H-2} \\
\asymp & n^{-1} \sum_{j=1}^{n-1} j^{2 H-2} \asymp n^{-1}+n^{2 H-2}
\end{aligned}
$$

(we used the fact that $r+1+\beta+\delta \geqslant 1$ since $r, \beta, \delta \geqslant 0$ ). For the fifth sum, we can proceed analogously and we also obtain $n^{-1}+n^{2 H-2}$ for bound.

For the sixth sum, we have

$$
\begin{aligned}
& n^{-2} \sum_{\substack{i=j \\
k=\ell \\
k \neq i}} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \\
\asymp & n^{-2} \sum_{k \neq i} \rho_{H}(k-i)^{2 q-2-2 r} \asymp n^{-2} \sum_{k \neq i}|k-i|^{(2 q-2-2 r)(2 H-2)} \asymp n^{-2} \sum_{k \neq i}|k-i|^{4 H-4} \\
\asymp & n^{-1} \sum_{k=1}^{n-1} k^{4 H-4} \asymp n^{-1}+n^{4 H-4}
\end{aligned}
$$

(here, we used $r \leqslant q-2$ ). For the seventh and the eighth sums, we can proceed analogously and we also obtain $n^{-1}+n^{4 H-4}$ for bound.

For the ninth sum, we have

$$
\begin{aligned}
& n^{-2} \sum_{\substack{i=j, k \neq i \\
k \neq \ell, \ell \neq i}} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \\
\asymp & n^{-2} \sum_{\substack{k \neq i \\
k \neq \ell, \ell \neq i}} \rho_{H}(k-\ell)^{r+1} \rho_{H}(k-i)^{q-r-1} \rho_{H}(\ell-i)^{q-r-1} .
\end{aligned}
$$

Now, let us decompose the sum $\sum_{k \neq i, k \neq \ell, \ell \neq i}$ into

$$
\sum_{k>\ell>i}+\sum_{k>i>\ell}+\sum_{\ell>i>k}+\sum_{\ell>k>i}+\sum_{i>\ell>k}+\sum_{i>k>\ell}
$$

For the first term (for instance), we have

$$
\begin{aligned}
& n^{-2} \sum_{k>\ell>i} \rho_{H}(k-\ell)^{r+1} \rho_{H}(k-i)^{q-r-1} \rho_{H}(\ell-i)^{q-r-1} \\
\asymp & n^{-2} \sum_{k>\ell>i}(k-\ell)^{(r+1)(2 H-2)}(k-i)^{(q-r-1)(2 H-2)}(\ell-i)^{(q-r-1)(2 H-2)} \\
\asymp & n^{-2} \sum_{k>\ell>i}(k-\ell)^{q(2 H-2)}(\ell-i)^{(q-r-1)(2 H-2)} \quad \text { since } k-i>k-\ell \\
= & n^{-2} \sum_{k} \sum_{\ell<k}(k-\ell)^{q(2 H-2)} \sum_{i<\ell}(\ell-i)^{(q-r-1)(2 H-2)} \\
\asymp & n^{-2} \sum_{k} \sum_{\ell<k}(k-\ell)^{q(2 H-2)} \sum_{i<\ell}(\ell-i)^{2 H-2} \quad \text { since } q-r-1 \geq 1 \\
\asymp & n^{-1} \sum_{\ell=1}^{n-1} \ell^{2 q H-2 q} \sum_{i=1}^{n-1} i^{2 H-2} \\
\asymp & n^{-1}\left(1+n^{2 q H-2 q+1}\right)\left(1+n^{2 H-1}\right) \asymp n^{-1}+n^{2 H-2} \quad \text { since } 2 q H-2 q+1<0 .
\end{aligned}
$$

We obtain the same bound for the other terms. By proceeding in the same way than for the ninth term, we also obtain the bound $n^{-1}+n^{2 H-2}$ for the tenth, eleventh, twelfth, thirteenth and fourteenth terms.
For the fifteenth (and last!) sum, we decompose $\sum_{(i, j, k, \ell}$ are all different) as follows

$$
\begin{equation*}
\sum_{k>l>i>j}+\sum_{k>l>j>i}+\ldots \tag{4.57}
\end{equation*}
$$

For the first term, we have:

$$
\begin{aligned}
& n^{-2} \sum_{k>\ell>i>j} \rho_{H}(k-\ell)^{r+1} \rho_{H}(i-j)^{r+1} \rho_{H}(k-i)^{\alpha} \rho_{H}(k-j)^{\beta} \rho_{H}(\ell-i)^{\gamma} \rho_{H}(\ell-j)^{\delta} \\
\asymp & n^{-2} \sum_{k>\ell>i>j}(k-\ell)^{q(2 H-2)}(i-j)^{(r+1)(2 H-2)}(\ell-i)^{(q-r-1)(2 H-2)} \\
= & n^{-2} \sum_{k} \sum_{\ell<k}(k-\ell)^{q(2 H-2)} \sum_{i<\ell}(\ell-i)^{(q-r-1)(2 H-2)} \sum_{j<i}(i-j)^{(r+1)(2 H-2)} \\
\asymp & n^{-1} \sum_{\ell=1}^{n-1} \ell^{q(2 H-2)} \sum_{i=1}^{n-1} i^{(q-r-1)(2 H-2)} \sum_{j=1}^{n-1} j^{(r+1)(2 H-2)} \\
\asymp & n^{-1}\left(1+n^{2 q H-2 q+1}\right)\left(1+n^{(q-r-1)(2 H-2)+1}\right)\left(1+n^{(r+1)(2 H-2)+1}\right) \\
\asymp & n^{-1}\left(1+n^{2 H-1}+n^{2 q H-2 q+2}\right) \quad \text { since } 2 q H-2 q+1<0 \text { and } r+1, q-r-1 \geqslant 1 \\
\asymp & n^{-1}+n^{2 H-2}+n^{2 q H-2 q+1} .
\end{aligned}
$$

The same bound also holds for the other terms in (4.57). By combining all these bounds, we obtain

$$
\max _{r=1, \ldots, q-1} E\left|A_{r}(n)\right|^{2} \asymp n^{-1}+n^{2 H-2}+n^{2 q H-2 q+1},
$$

that finally gives:

$$
E\left(\frac{1}{q}\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}-1\right)^{2} \asymp n^{-1}+n^{2 H-2}+n^{2 q H-2 q+1} .
$$

The proof of Theorem 4.1 is now completed by means of Proposition 3.2.

## 5 Two proofs

### 5.1 Proof of Lemma 1.3

Proof of Point (i). Observe first that, for every $\nu>0$, the random variable $F^{*}(\nu):=F(\nu)+\nu$ has a non-centered Gamma law with parameter $\nu / 2$. The fact that

$$
E\left[2 F^{*}(\nu) f^{\prime}\left(F^{*}(\nu)-\nu\right)\right]=E\left[2\left(F^{*}(\nu)-\nu\right) f^{\prime}\left(F^{*}(\nu)\right)\right],
$$

for every $f$ as in the statement, is therefore an immediate consequence of [45, Proposition 1 and Section 4(2)]. Now suppose that $W$ verifies (1.13). By choosing $f$ with support in $(-\infty,-\nu)$, one deduces immediately that $P(W \leqslant-\nu)=0$. To conclude, we apply once again the results contained in [45], to infer that the relations

$$
P(W \leqslant-\nu)=0 \quad \text { and } \quad E\left[2(W+\nu) f^{\prime}(W)-W f(W)\right]=0
$$

imply that, necessarily, $W+\nu \stackrel{\text { Law }}{=} F^{*}(\nu)$.
Proof of Point (ii). Fix $\nu>0$, consider a function $h$ as in the statement and define $h_{\nu}(y)=$ $\overline{h(y-\nu), y>0}$. Plainly, $h_{\nu}$ is twice differentiable, and $\left|h_{\nu}(y)\right| \leqslant c \exp \{-\nu a\} \exp \{a y\}, y>0$ (recall that $a>1 / 2$ ). In view of these properties, according to Luk [24, Th. 1], the secondorder Stein equation

$$
\begin{equation*}
h_{\nu}(y)-E\left(h_{\nu}\left(F^{*}(\nu)\right)=2 y g^{\prime \prime}(y)-(y-\nu) g^{\prime}(y), \quad y>0,\right. \tag{5.58}
\end{equation*}
$$

(where, as before, we set $F^{*}(\nu)=F(\nu)+\nu$ ) admits a solution $g$ such that $\left\|g^{\prime}\right\|_{\infty} \leqslant 2\left\|h^{\prime}\right\|_{\infty}$ and $\left\|g^{\prime \prime}\right\|_{\infty} \leqslant\left\|h^{\prime \prime}\right\|_{\infty}$. Since $f(x)=g^{\prime}(x+\nu), x>-\nu$, is a solution of (1.12), the conclusion is immediately obtained.
Proof of Point (iii). According to a result of Pickett [40], as reported in [42, Lemma 3.1], when $\nu \geqslant 1$ is an integer, the ancillary Stein equation (5.58) admits a solution $g$ such that $\left\|g^{\prime}\right\|_{\infty} \leqslant \sqrt{2 \pi / \nu}\|h\|_{\infty}$ and $\left\|g^{\prime \prime}\right\|_{\infty} \leqslant \sqrt{2 \pi / \nu}\left\|h^{\prime}\right\|_{\infty}$. The conclusion is obtained as in the proof of Point (ii).

### 5.2 Proof of Theorem 1.5

We begin with a technical lemma.
Lemma 5.1 Let $F=I_{2}(f)$ be a random variable living in the second Wiener chaos of an isonormal Gaussian process $X$ (over a real Hilbert space $\mathfrak{H}$ ). Then

$$
\begin{equation*}
E\left(\|D F\|_{\mathfrak{H}}^{4}\right)=\frac{2}{3} E\left(F^{4}\right)+2 E\left(F^{2}\right)^{2} . \tag{5.59}
\end{equation*}
$$

Proof. Without loss of generality, we can assume that $\mathfrak{H}=\mathrm{L}^{2}(A, \mathscr{A}, \mu)$, where $(A, \mathscr{A})$ is a measurable space, and $\mu$ is a $\sigma$-finite and non-atomic measure. On one hand, thanks to the multiplication formula (2.29), we can write

$$
F^{2}=I_{4}(f \otimes f)+4 I_{2}\left(f \otimes_{1} f\right)+E\left(F^{2}\right) .
$$

In particular, this yields

$$
L\left(F^{2}\right)=-4 I_{4}(f \otimes f)-8 I_{2}\left(f \otimes_{1} f\right) .
$$

On the other hand, (2.32) implies that $D_{a} F=2 I_{1}(f(\cdot, a))$. Consequently, again by (2.29):

$$
\begin{align*}
\|D F\|_{\mathfrak{H}}^{2} & =4 \int_{A} I_{1}(f(\cdot, a))^{2} \mu(d a) \\
& =4 \int_{A} I_{2}(f(\cdot, a) \otimes f(\cdot, a)) \mu(d a)+E\left(\|D F\|_{\mathfrak{H}}^{2}\right) \\
& =4 I_{2}\left(f \otimes_{1} f\right)+2 E\left(F^{2}\right), \quad \text { by }(2.34) \text { and since } \int_{A} f(\cdot, a) \otimes f(\cdot, a) \mu(d a)=f \otimes_{1} f . \tag{5.60}
\end{align*}
$$

Taking into account the orthogonality between multiple stochastic integrals of different orders, we deduce

$$
\begin{equation*}
E\left[\|D F\|_{\mathfrak{H}}^{2} L\left(F^{2}\right)\right]=-32 E\left[\left(I_{2}\left(f \otimes_{1} f\right)\right)^{2}\right]=-2 E\left[\|D F\|_{\mathfrak{S}}^{2}\left(F^{2}-E\left(F^{2}\right)\right)\right] . \tag{5.61}
\end{equation*}
$$

Finally, we have

$$
\begin{aligned}
E\left[\|D F\|_{\mathfrak{H}}^{4}\right] & =E\left[\|D F\|_{\mathfrak{H}}^{2}\langle D F, D F\rangle_{\mathfrak{H}}\right] \\
& =E\left[\|D F\|_{\mathfrak{H}}^{2}\left(\delta D F \times F-\frac{1}{2} \delta D\left(F^{2}\right)\right)\right] \quad \text { by identity (2.31), } \\
& =2 E\left[\|D F\|_{\mathfrak{H}}^{2} F^{2}\right]+\frac{1}{2} E\left[\|D F\|_{\mathfrak{H}}^{2} L\left(F^{2}\right)\right] \quad \text { using } \delta D=-L, \\
& =E\left[\|D F\|_{\mathfrak{H}}^{2} F^{2}\right]+E\left(F^{2}\right) E\left[\|D F\|_{\mathfrak{H}}^{2}\right] \quad \text { using (5.61), } \\
& =\frac{2}{3} E\left(F^{4}\right)+2 E\left(F^{2}\right)^{2} \quad \text { by (2.34). }
\end{aligned}
$$

Now, let us go back to the proof of the first point in Theorem 1.5. In view of Theorem 3.1, it is sufficient to prove that

$$
\begin{equation*}
E\left(\left|1-\frac{1}{2}\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}\right|^{2}\right) \leqslant \frac{1}{6}\left|E\left(Z_{n}^{4}\right)-3\right|+\frac{3+E\left(Z_{n}^{2}\right)}{2}\left|E\left(Z_{n}^{2}\right)-1\right| . \tag{5.62}
\end{equation*}
$$

We have

$$
\begin{aligned}
& E\left(\left|1-\frac{1}{2}\left\|D Z_{n}\right\|_{\mathfrak{F}}^{2}\right|^{2}\right) \\
& =1-E\left(\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}\right)+\frac{1}{4} E\left(\left\|D Z_{n}\right\|_{\mathfrak{H}}^{4}\right) \\
& =1-2 E\left(Z_{n}^{2}\right)+\frac{1}{6} E\left(Z_{n}^{4}\right)+\frac{1}{2} E\left(Z_{n}^{2}\right)^{2} \quad \text { by }(2.34) \text { and (5.59) } \\
& =\frac{1}{6}\left(E\left(Z_{n}^{4}\right)-3\right)+\left(E\left(Z_{n}^{2}\right)-1\right)\left(\frac{1}{2} E\left(Z_{n}^{2}\right)-\frac{3}{2}\right) .
\end{aligned}
$$

The estimate (5.62) follows immediately.
Similarly, for the second point of Theorem 1.5, it is sufficient to prove (see Proposition 3.13) that

$$
\begin{align*}
E\left(\mid 2 Z_{n}-2 \nu-\right. & \left.\left.\frac{1}{2}\left\|D Z_{n}\right\|_{\mathfrak{H}}\right|^{2}\right)  \tag{5.63}\\
& \leqslant \frac{1}{6}\left|E\left(Z_{n}^{4}\right)-12 E\left(Z_{n}^{3}\right)-12 \nu^{2}+48 \nu\right|+\frac{\left|8-6 \nu+E\left(Z_{n}^{2}\right)\right|}{2}\left|E\left(Z_{n}^{2}\right)-2 \nu\right|
\end{align*}
$$

By using the relations

$$
\begin{aligned}
& E\left(\left|2 Z_{n}-2 \nu-\frac{1}{2}\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}\right|^{2}\right) \\
= & 4 E\left(Z_{n}^{2}\right)+4 \nu^{2}+\frac{1}{4} E\left(\left\|D Z_{n}\right\|_{\mathfrak{H}}^{4}\right)-2 E\left(Z_{n}\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}\right)-2 \nu E\left(\left\|D Z_{n}\right\|_{\mathfrak{H}}^{2}\right) \\
= & 4(1-\nu) E\left(Z_{n}^{2}\right)+4 \nu^{2}+\frac{1}{6} E\left(Z_{n}^{4}\right)+\frac{1}{2} E\left(Z_{n}^{2}\right)^{2}-2 E\left(Z_{n}^{3}\right) \quad \text { by }(2.34) \text { and }(5.59) \\
= & \left(E\left(Z_{n}^{2}\right)-2 \nu\right)\left(4-3 \nu+\frac{1}{2} E\left(Z_{n}^{2}\right)\right)+\frac{1}{6}\left(E\left(Z_{n}^{4}\right)-12 E\left(Z_{n}^{3}\right)-12 \nu^{2}+48 \nu\right),
\end{aligned}
$$

the estimate (5.63) follows immediately.

Acknowledgment: We are grateful to G. Reinert for bringing to our attention references [24] and [40]. We also thank J. Dedecker for useful discussions.
This paper is dedicated to the memory of Livio Zerbini.

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