Malliavin-Stein Method:
a Survey of Recent Developments

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Abstract

Initiated around the year 2007, the Malliavin-Stein approach to probabilistic approximations combines Stein’s method with infinite-dimensional integration by parts formulae based on the use of Malliavin-type operators. In the last decade, Malliavin-Stein techniques have allowed researchers to establish new quantitative limit theorems in a variety of domains of theoretical and applied stochastic analysis. The aim of this short survey is to illustrate some of the latest developments of the Malliavin-Stein method, with specific emphasis on extensions and generalisations in the framework of Markov semigroups and of random point measures.

Keywords: Limit Theorems; Stein’s method; Malliavin Calculus; Wiener Space; Poisson Space; Multiple Integral; Markov Triple; Markov Generator; Eigenspace; Eigenfunction; Spectrum; Functional Γ-Calculus; Weak Convergence; Fourth Moment Theorems; Berry–Essen Bounds; Probability Metrics.

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1 Introduction and overview

The Malliavin-Stein method for probabilistic approximations was initiated in the paper [NP09b], with the aim of providing a quantitative counterpart to the (one- and multi-dimensional) central limit theorems for random variables living in the Wiener chaos of a general separable Gaussian field. As formally discussed in the sections to follow, the basic idea of the approach initiated in [NP09b] is that, in order to assess the discrepancy between some target law (Normal or Gamma, for instance), and the distribution of a non-linear functional of a Gaussian field, one can fruitfully apply infinite-dimensional integration by parts formulae from the Malliavin calculus of variations [M97, NP12, Nua06, Nua09] to the general bounds associated with the so-called Stein’s method for probabilistic approximations [NP12, CGS10]. In particular, the Malliavin-Stein approach captures and amplifies the essence of [C09], where Stein’s method was combined with finite-dimensional integration by parts formulae for Gaussian vectors, in order to deduce second order Poincaré inequalities – as applied to random matrix models with Gaussian-subordinate entries (see also [Vid17]).

We recall that, as initiated by P. Malliavin in the path-breaking reference [M78], the Malliavin calculus is an infinite-dimensional differential calculus, whose operators act on smooth non-linear functionals of Gaussian fields (or of more general probabilistic objects). As vividly described in the classical references [M97, Nua06], as well as in the more recent books [NP12, Nua09], since its inception such a theory has generated an astonishing number of applications, ranging e.g. from mathematical physics to stochastic differential equations, and from mathematical finance to stochastic geometry, analysis on manifolds and mathematical statistics. On the other hand, the similarly successful and popular Stein’s method (as created by Ch. Stein in the classical reference [S72] – see also the 1986 monograph [S86]) is a collection of analytical techniques, allowing one to estimate the distance between the distributions of two random objects, by using characterising differential operators. The discovery in [NP09b] that the two theories can be fruitfully combined has been a major breakthrough in the domain of probabilistic limit theorems and approximations.

Since the publication of [NP09b], the Malliavin-Stein method has generated several hundreds of papers, with ramifications in many (often unexpected) directions, including functional inequalities, random matrix theory, stochastic geometry, non-commutative probability and computer sciences. These developments largely exceed the scope of the present survey, and we invite the interested reader to consult the following references (i)–(iii) for a more detailed presentation: (i) the webpage [WWW] is a constantly updated resource, listing all existing papers written around the Malliavin-Stein method; (ii) the monograph [NP12], written in 2012, contains a self-contained presentation of Malliavin calculus and Stein’s method, as applied to functionals of general Gaussian fields, with specific emphasis on random variables belonging to a fixed Wiener chaos; (iii) the text [PR16] is a collection of surveys, containing an in-depth presentation of variational techniques on the Poisson spaces (including the Malliavin-Stein method), together with their application to asymptotic problems arising in stochastic geometry.

The aim of the present survey is twofolds. On the one hand, we aim at presenting the essence of the Malliavin-Stein’s method for functionals of Gaussian fields, by discussing the crucial elements of Malliavin calculus and Stein’s method together with their interaction (see Section 2 and Section 3). On the other hand, we aim at introducing the reader to some of the most recent developments on the theory, with specific focus on the general theory of Markov semigroups in a diffusive setting (following the seminal references [Led12, ACP14], as well as [NPS15, LNP15, LNP16]), and on integration by parts formulae (and associated operators) in the context of functionals of a random point measure [DP18b, DVZ18]. This corresponds to the content of Section 4 and Section 5, respectively. Finally, Section 5 deals with some new results (and open problems) concerning $\chi^2$ approximations.
From now on, every random object will be defined on a suitable common probability space \((\Omega, \mathcal{F}, P)\), with \(E\) indicating mathematical expectation with respect to \(P\). Throughout the paper, the symbol \(\mathcal{N}(\mu, \sigma^2)\) will be shorthand for the one-dimensional Gaussian distribution with mean \(\mu \in \mathbb{R}\) and variance \(\sigma^2 > 0\). In particular, \(X \sim \mathcal{N}(\mu, \sigma^2)\) if and only if
\[
P[X \in A] = \int_A e^{-\frac{(x-\mu)^2}{2\sigma^2}} \frac{dx}{\sqrt{2\pi\sigma^2}},
\]
for every Borel set \(A \subset \mathbb{R}\).

2 Elements of Stein’s method for normal approximations

In this section, we briefly introduce the main ingredients of Stein’s method for normal approximations in dimension one. The approximation will be performed with respect to the total variation and 1-Wasserstein distances between the distributions of two random variables; more detailed informations about these distances can be found in [NP12, Appendix C] and the references therein.

The crucial intuition behind Stein’s method lies in the following heuristic reasoning: it is a well-known fact (see e.g. Lemma 2.1-(e) below) that a random variable \(X\) has the standard \(\mathcal{N}(0,1)\) distribution if and only if
\[
\mathbb{E}[X f(X) - f'(X)] = 0,
\]
for every smooth mapping \(f : \mathbb{R} \to \mathbb{R}\); heuristically, it follows that, if \(X\) is a random variable such that the quantity \(\mathbb{E}[X f(X) - f'(X)]\) is close to zero for a large class of test functions \(f\), then the distribution of \(X\) should be close to Gaussian.

The fact that such a heuristic argument can be made rigorous and applied in a wide array of probabilistic models was the main discovery of Stein’s original contribution [S72], where the foundations of Stein’s method were first laid. The reader is referred to Stein’s monograph [S86], as well as the recent books [CGS10, NP12], for an exhaustive presentation of the theory and its applications (in particular, for extensions to multidimensional approximations).

We recall that the total variation distance, between the laws of two real-valued random variables \(F\) and \(G\), is defined by
\[
d_{TV}(F, G) := \sup_{B \in \mathcal{B}(\mathbb{R})} \left| \mathbb{P}(F \in B) - \mathbb{P}(G \in B) \right|.
\]
One has to note that the topology induced by the distance \(d_{TV}\) – on the set of all probability measures on \(\mathbb{R}\) – is stronger than the topology of convergence in distribution; we will often use the following equivalent representation of \(d_{TV}\) (see e.g. [NP12, p. 213]):
\[
d_{TV}(F, G) = \frac{1}{2} \sup \left\{ \left| \mathbb{E}[h(F)] - \mathbb{E}[h(G)] \right| : h \text{ is Borel measurable and } \|h\|_{\infty} \leq 1 \right\}.
\]

The 1-Wasserstein distance \(d_W\), between the distributions of two real-valued integrable random variables \(F\) and \(G\), is given by
\[
d_W(F, G) := \sup_{h \in \text{Lip}(1)} \left| \mathbb{E}[h(F)] - \mathbb{E}[h(G)] \right|,
\]
where \(\text{Lip}(K), K > 0\) stands for the class of all Lipschitz mappings \(h : \mathbb{R} \to \mathbb{R}\) such that \(h\) has a Lipschitz constant \(\leq K\). As for total variation, the topology induced by \(d_W\) – on the set of all probability measures on \(\mathbb{R}\) having a finite absolute first moment – is stronger than the topology of convergence in distribution; it is also interesting to recall the dual representation
\[
d_W(F, G) = \inf \mathbb{E} |X - Y|,
\]
where the infimum is taken over all couplings \((X, Y)\) of \(F\) and \(G\); see e.g. [Vil09, p. 95] for a discussion of this fact.

The following classical result, whose complete proof can be found e.g. in [NP12, p. 64 and p. 67], contains all the elements of Stein’s method that are needed for our discussion; as for many fundamental findings in the area, such a result can be traced back to [S72].

**Lemma 2.1.** Let \(N \sim \mathcal{N}(0, 1)\) be a standard Gaussian random variable.

(a) Fix \(h : \mathbb{R} \to [0, 1]\) a Borel-measurable function. Define \(f_h : \mathbb{R} \to \mathbb{R}\) as

\[
f_h(x) := e^{-\frac{x^2}{2}} \int_{-\infty}^{x} \{h(y) - \mathbb{E}[h(N)]\} e^{-\frac{y^2}{2}} dy, \quad x \in \mathbb{R}.
\]

(2.6)

Then, \(f_h\) is continuous on \(\mathbb{R}\) with \(\|f_h\|_{\infty} \leq \sqrt{\frac{\pi}{2}}\) and \(f_h \in \text{Lip}(2)\). Moreover, there exists a version of \(f_h'\) verifying

\[
f_h'(x) - xf_h(x) = h(x) - \mathbb{E}[h(N)], \quad \text{for all } x \in \mathbb{R}.
\]

(2.7)

(b) Consider \(h : \mathbb{R} \to \mathbb{R} \in \text{Lip}(1)\), and define \(f_h : \mathbb{R} \to \mathbb{R}\) as in (2.6). Then, \(f_h\) is of class \(C^1\) on \(\mathbb{R}\), with \(\|f_h'\|_{\infty} \leq 1\) and \(f_h' \in \text{Lip}(2)\), and \(f_h\) solves (2.7).

(c) Let \(X\) be an integrable random variable. Then

\[
d_{TV}(X, N) \leq \sup_{f} \left| \mathbb{E}[f(X)X - f'(X)] \right|
\]

where the supremum is taken over all pairs \((f, f')\) such that \(f\) is a Lipschitz function whose absolute value is bounded by \(\sqrt{\frac{\pi}{2}}\), and \(f'\) is a version of the derivative of \(f\) satisfying \(\|f'\| \leq 2\).

(d) Let \(X\) be an integrable random variable. Then,

\[
d_W(X, N) \leq \sup_{f} \left| \mathbb{E}[f(X)X - f'(X)] \right|
\]

where the supremum is taken over all \(C^1\) functions \(f : \mathbb{R} \to \mathbb{R}\) such that \(\|f'\| \leq 2\) and \(f' \in \text{Lip}(2)\).

(e) Let \(X\) be a general random variable. Then \(X \sim \mathcal{N}(0, 1)\) if and only if \(\mathbb{E}[f'(X) - Xf(X)] = 0\) for every absolutely continuous function \(f\) such that \(\mathbb{E}|f'(N)| < +\infty\).

**Sketch of the proof.** Points (a) and (b) can be verified by a direct computation. Point (c) and Point (d) follow by plugging the left-hand side of (2.7) into (2.3) and (2.4), respectively. Finally, the fact that the relation \(\mathbb{E}[f'(X) - Xf(X)] = 0\) implies that \(X \sim \mathcal{N}(0, 1)\) is a direct consequence of Point (c), whereas the reverse implication follows by an integration by parts argument.

### 3 Normal approximation with Stein’s method and Malliavin Calculus

The first part of the present section contains some elements of Gaussian analysis and Malliavin calculus. The reader can consult for instance the references [NP12, Nua06, M97, Nua09] for further details. In Section 3.2, we will shortly explore the connection between Malliavin calculus and the version Stein’s method presented in Section 2.
3.1 Isonormal processes, multiple integrals, and the Malliavin operators

Let \( \mathcal{H} \) be a real separable Hilbert space. For any \( q \geq 1 \), we write \( \mathcal{H}^{\otimes q} \) and \( \mathcal{H}^{\ominus q} \) to indicate, respectively, the \( q \)th tensor power and the \( q \)th symmetric tensor power of \( \mathcal{H} \); we also set by convention \( \mathcal{H}^{\otimes 0} = \mathcal{H}^{\ominus 0} = \mathbb{R} \). When \( \mathcal{H} = L^2(\mathcal{A}, \mu) := L^2(\mu) \), where \( \mu \) is a \( \sigma \)-finite and non-atomic measure on the measurable space \( (\mathcal{A}, \mathcal{A}) \), then \( \mathcal{H}^{\otimes q} \approx L^2(A^q, \mathcal{A}^q, \mu^q) := L^2_0(\mu^q) \), and \( \mathcal{H}^{\ominus q} \approx L^2(A^q, \mathcal{A}^q, \mu^q) := L^2_0(\mu^q) \), where \( L^2_0(\mu^q) \) stands for the subspace of \( L^2(\mu^q) \) composed of those functions that are \( \mu^q \)-almost everywhere symmetric. We denote by \( W = \{W(h) : h \in \mathcal{H}\} \) an isonormal Gaussian process over \( \mathcal{H} \). This means that \( W \) is a centered Gaussian family with a covariance structure given by the relation \( \mathbb{E}[W(h)W(g)] = \langle h, g \rangle_\mathcal{H} \). Without loss of generality, we can also assume that \( \mathcal{F} = \sigma(W) \), that is, \( \mathcal{F} \) is generated by \( W \), and use the shorthand notation \( L^2(\Omega) := L^2(\Omega, \mathcal{F}, \mathbb{P}) \).

For every \( q \geq 1 \), the symbol \( C_q \) stands for the \( q \)th Wiener chaos of \( W \), defined as the closed linear subspace of \( L^2(\Omega) \) generated by the family \( \{H_q(W(h)) : h \in \mathcal{H}, \|h\|_\mathcal{H} = 1\} \), where \( H_q \) is the \( q \)th Hermite polynomial, defined as follows:

\[
H_q(x) = (-1)^q e^{\frac{x^2}{2}} \frac{d^q}{dx^q} (e^{-\frac{x^2}{2}}).
\]

We write by convention \( C_0 = \mathbb{R} \). For any \( q \geq 1 \), the mapping \( I_q(h^{\otimes q}) = H_q(W(h)) \) can be extended to a linear isometry between the symmetric tensor product \( \mathcal{H}^{\otimes q} \) equipped with the modified norm \( \sqrt{q!} \|\cdot\|_{\mathcal{H}^{\otimes q}} \) and the \( q \)th Wiener chaos \( C_q \). For \( q = 0 \), we write by convention \( I_0(c) = c, c \in \mathbb{R} \).

It is well-known that \( L^2(\Omega) \) can be decomposed into the infinite orthogonal sum of the spaces \( C_q \); this means that any square-integrable random variable \( F \in L^2(\Omega) \) admits the following Wiener-Itô chaotic expansion

\[
F = \sum_{q=0}^{\infty} I_q(f_q),
\]

where the series converges in \( L^2(\Omega) \), \( f_0 = \mathbb{E}[F] \), and the kernels \( f_q \in \mathcal{H}^{\otimes q}, q \geq 1 \), are uniquely determined by \( F \). For every \( q \geq 0 \), we denote by \( J_q \) the orthogonal projection operator on the \( q \)th Wiener chaos. In particular, if \( F \in L^2(\Omega) \) has the form \( \sum_{q=0}^{\infty} I_q(f_q) \) for every \( q \geq 0 \).

Let \( \{e_k, k \geq 1\} \) be a complete orthonormal system in \( \mathcal{H} \). Given \( f \in \mathcal{H}^{\otimes p} \) and \( g \in \mathcal{H}^{\otimes q} \), for every \( r = 0, \ldots, p \wedge q \), the contraction of \( f \) and \( g \) of order \( r \) is the element of \( \mathcal{H}^{\otimes (p+q-2r)} \) defined by

\[
f \otimes_r g = \sum_{i_1, \ldots, i_r = 1}^{\infty} \langle f, e_{i_1} \otimes \cdots \otimes e_{i_r} \rangle_{\mathcal{H}^{\otimes r}} \otimes \langle g, e_{i_1} \otimes \cdots \otimes e_{i_r} \rangle_{\mathcal{H}^{\otimes r}}.
\]

Notice that the definition of \( f \otimes_r g \) does not depend on the particular choice of \( \{e_k, k \geq 1\} \), and that \( f \otimes_r g \) is not necessarily symmetric; we denote its symmetrization by \( f \otimes_r g \in \mathcal{H}^{\otimes (p+q-2r)} \). Moreover, \( f \otimes_0 g = f \otimes g \) equals the tensor product of \( f \) and \( g \) while, for \( p = q \), \( f \otimes_q g = \langle f, g \rangle_{\mathcal{H}^{\otimes q}} \). When \( \mathcal{H} = L^2(A, \mathcal{A}, \mu) \) and \( r = 1, \ldots, p \wedge q \), the contraction \( f \otimes_r g \) is the element of \( L^2(\mu^{p+q-2r}) \) given by

\[
f \otimes_r g(x_1, \ldots, x_{p+q-2r}) = \int_{A^r} f(x_1, \ldots, x_{p-r}, a_1, \ldots, a_r) \times g(x_{p-r+1}, \ldots, x_{p+q-2r}, a_1, \ldots, a_r) d\mu(a_1) \ldots d\mu(a_r).
\]

It is a standard fact of Gaussian analysis that the following multiplication formula holds: if \( f \in \mathcal{H}^{\otimes p} \) and \( g \in \mathcal{H}^{\otimes q} \), then

\[
I_p(f)I_q(g) = \sum_{r=0}^{p \wedge q} r! \binom{p}{r} \binom{q}{r} I_{p+q-2r}(f \otimes_r g).
\]
We now introduce some basic elements of the Malliavin calculus with respect to the isonormal Gaussian process \( W \).

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

\[
F = g(W(\varphi_1), \ldots, W(\varphi_n)),
\]

where \( n \geq 1, \ g : \mathbb{R}^n \to \mathbb{R} \) is an infinitely differentiable function such that its partial derivatives have polynomial growth, and \( \varphi_i \in \mathcal{H}, \ i = 1, \ldots, n \). The Malliavin derivative of \( F \) with respect to \( W \) is the element of \( L^2(\Omega) \) defined as

\[
DF = \sum_{i=1}^n \frac{\partial g}{\partial x_i} (W(\varphi_1), \ldots, W(\varphi_n)) \varphi_i.
\]

In particular, \( DW(h) = h \) for every \( h \in \mathcal{H} \). By iteration, one can define the \( m \)th derivative \( D^m F \), which is an element of \( L^2(\Omega) \otimes^m \), for every \( m \geq 2 \). For \( m \geq 1 \) and \( p \geq 1 \), \( \mathbb{D}^{m,p} \) denotes the closure of \( \mathcal{S} \) with respect to the norm \( \| \cdot \|_{m,p} \), defined by the relation

\[
\|F\|_{m,p}^p = \mathbb{E}[|F|^p] + \sum_{i=1}^m \mathbb{E}[\|D^i F\|_{\mathcal{H}^{\otimes i}}^p].
\]

We often use the (canonical) notation \( \mathbb{D}^\infty := \bigcap_{m=1}^\infty \bigcap_{p=1}^{\infty} \mathbb{D}^{m,p} \). For example, it is a well-known fact that any random variable \( F \) that is a finite linear combination of multiple Wiener-Itô integrals is an element of \( \mathbb{D}^\infty \). The Malliavin derivative \( D \) obeys the following chain rule. If \( \phi : \mathbb{R}^n \to \mathbb{R} \) is continuously differentiable with bounded partial derivatives and if \( F = (F_1, \ldots, F_n) \) is a vector of elements of \( \mathbb{D}^{1,2} \), then \( \phi(F) \in \mathbb{D}^{1,2} \) and

\[
D \phi(F) = \sum_{i=1}^n \frac{\partial \phi}{\partial x_i}(F)DF_i.
\]

Note also that a random variable \( F \) as in (3.6) is in \( \mathbb{D}^{1,2} \) if and only if \( \sum_{q=1}^\infty q\|J_q F\|_{L^2(\Omega)}^2 < \infty \) and in this case one has the following explicit relation:

\[
\mathbb{E}[\|DF\|_{\mathcal{H}}^2] = \sum_{q=1}^\infty q\|J_q F\|_{L^2(\Omega)}^2.
\]

If \( \mathcal{H} = L^2(A, \mathcal{A}, \mu) \) (with \( \mu \) non-atomic), then the derivative of a random variable \( F \) as in (3.2) can be identified with the element of \( L^2(A \times \Omega) \) given by

\[
D_t F = \sum_{q=1}^\infty q I_{q-1}(f_q(\cdot,t)), \quad t \in A.
\]

The operator \( \mathbf{L} \), defined as \( \mathbf{L} = \sum_{q=0}^\infty -q J_q \), is the infinitesimal generator of the Ornstein-Uhlenbeck semigroup. The domain of \( \mathbf{L} \) is

\[
\text{Dom}\mathbf{L} = \{ F \in L^2(\Omega) : \sum_{q=1}^\infty q^2 \|J_q F\|_{L^2(\Omega)}^2 < \infty \} = \mathbb{D}^{2,2}.
\]

For any \( F \in L^2(\Omega) \), we define \( \mathbf{L}^{-1} F = \sum_{q=1}^\infty -\frac{1}{q} J_q(F) \). The operator \( \mathbf{L}^{-1} \) is called the pseudo-inverse of \( \mathbf{L} \). Indeed, for any \( F \in L^2(\Omega) \), we have that \( \mathbf{L}^{-1} F \in \text{Dom}\mathbf{L} = \mathbb{D}^{2,2} \), and

\[
\mathbf{L}\mathbf{L}^{-1} F = F - \mathbb{E}(F).
\]

The following infinite dimensional Malliavin integration by parts formula plays a crucial role in the analysis (see for instance [NP12, Section 2.9] for a proof).
Lemma 3.1. Suppose that $F \in \mathbb{D}^{1,2}$ and $G \in L^2(\Omega)$. Then, $L^{-1}G \in \mathbb{D}^{2,2}$ and

$$
\mathbb{E}[FG] = \mathbb{E}[F] \mathbb{E}[G] + \mathbb{E}[(DF, -DL^{-1}G)_\mathbb{B}].
$$

(3.10)

Inspired by the Malliavin integration by parts formula appearing in Lemma 3.1, we now introduce a class of iterated Gamma operators. We will need such operators in Section 6.

Definition 3.2 (See Chapter 8 in [NP12]). Let $F \in \mathbb{D}^{\infty}$; the sequence of random variables $\{\Gamma_i(F)\}_{i \geq 0} \subset \mathbb{D}^{\infty}$ is recursively defined as follows. Set $\Gamma_0(F) = F$ and, for every $i \geq 1$,

$$
\Gamma_i(F) = \langle DF, -DL^{-1}\Gamma_{i-1}(F)\rangle_\mathbb{B}.
$$

Definition 3.3 (Cumulants). Let $F$ be a real-valued random variable such that $\mathbb{E}|F|^m < \infty$ for some integer $m \geq 1$, and write $\varphi_F(t) = \mathbb{E}[e^{itF}]$, $t \in \mathbb{R}$, for the characteristic function of $F$. Then, for $r = 1, \ldots, m$, the $r$th cumulant of $F$, denoted by $\kappa_r(F)$, is given by

$$
\kappa_r(F) = (-i)^r \frac{d^r}{dt^r} \log \varphi_F(t)|_{t=0}.
$$

(3.11)

Remark 3.4. When $\mathbb{E}(F) = 0$, then the first four cumulants of $F$ are the following: $\kappa_1(F) = \mathbb{E}[F] = 0$, $\kappa_2(F) = \mathbb{E}[F^2] = \text{Var}(F)$, $\kappa_3(F) = \mathbb{E}[F^3]$, and

$$
\kappa_4(F) = \mathbb{E}[F^4] - 3\mathbb{E}[F^2]^2.
$$

The following statement explicitly connects the expectation of the random variables $\Gamma_r(F)$ to the cumulants of $F$.

Proposition 3.5 (See Chapter 8 in [NP12]). Let $F \in \mathbb{D}^{\infty}$. Then $\kappa_r(F) = (r-1)!\mathbb{E}[\Gamma_{r-1}(F)]$ for every $r \geq 1$.

As announced, in the next subsection we show how to use the above Malliavin machinery in order to study the Stein’s bounds presented in Section 2.

3.2 Connection with Stein’s method

Let $F \in \mathbb{D}^{1,2}$ with $\mathbb{E}[F] = 0$ and $\mathbb{E}[F^2] = 1$. Take a $C^1$ function such that $\|f\| \leq \sqrt{2}$ and $\|f'\| \leq 2$. Using the Malliavin integration by parts formula stated in Lemma 3.1 together with the chain rule (3.7), we can write

$$
\left| \mathbb{E}[f'(F) - Ff(F)] \right| = \left| \mathbb{E}[f'(F) (1 - \langle DF, -DL^{-1}F\rangle_\mathbb{B})] \right|
\leq 2 \mathbb{E}\left| 1 - \langle DF, -DL^{-1}F\rangle_\mathbb{B} \right|.
$$

If we furthermore assume that $F \in \mathbb{D}^{1,4}$, then the random variable $1 - \langle DF, -DL^{-1}F\rangle_\mathbb{B}$ is square-integrable, using the Cauchy-Schwarz inequality we infer that

$$
\left| \mathbb{E}[f'(F) - Ff(F)] \right| \leq 2 \sqrt{\text{Var}(\langle DF, -DL^{-1}F\rangle_\mathbb{B})}.
$$

Note that in above we used the fact that $\mathbb{E}[\langle DF, -DL^{-1}F\rangle_\mathbb{B}] = \mathbb{E}[F^2] = 1$. The above arguments combined with Lemma 2.1 yield immediately the next crucial statement, originally proved in [NP09b].

1This is not completely accurate: attention has indeed to be paid to the fact that the function $f_h$ in (3.7) is only almost everywhere differentiable, and $F$ does not necessarily have a density – see [N12] Theorem 5.2 for a detailed proof based on Lusin Theorem.
Theorem 3.6. Let $F \in \mathbb{D}^{1,2}$ be a generic random element with $\mathbb{E}[F] = 0$ and $\mathbb{E}[F^2] = 1$. Let $N \sim \mathcal{N}(0, 1)$. Assume further that $F$ has a density with respect to the Lebesgue measure. Then,

$$d_{TV}(F, N) \leq 2 \mathbb{E}\left|1 - \langle DF, -DL^{-1}F\rangle_{\mathcal{S}}\right|.$$ 

Moreover, assume that $F \in \mathbb{D}^{1,4}$, then

$$d_{TV}(F, N) \leq 2\sqrt{\text{Var}(\langle DF, -DL^{-1}F\rangle_{\mathcal{S}})}.$$ 

In particular case, if $F = I_q(f)$ belongs to the Wiener chaos of order $q \geq 2$, then

$$d_{TV}(F, N) \leq 2\sqrt{\frac{q - 1}{3q} (\mathbb{E}[F^4] - 3)}.$$ 

(3.12)

Note that, by virtue of Lemma 2.1, similar bounds can be immediately obtained for the Wasserstein distance $d_W$ (and many more – see [NP12, Chapter 5]). In particular, the previous statement allows one to recover the following central limit theorem for chaotic random variables, first proved in [NP05].

Corollary 3.7 (Fourth Moment Theorem). Let $\{F_n\}_{n \geq 1} = \{I_q(f_n)\}_{n \geq 1}$ be a sequence of random elements in a fixed Wiener chaos of order $q \geq 2$ such that $\mathbb{E}[F_n^2] = q! \|f_n\|^2 = 1$. Assume that $N \sim \mathcal{N}(0, 1)$. Then, as $n$ tends to infinity, the following assertions are equivalent.

(I) $F_n \to N$ in distribution.

(II) $\mathbb{E}[F_n^4] \to 3 = \mathbb{E}[N^4]$.

As demonstrated by the webpage [WWW], the ‘fourth moment theorem’ stated in Corollary 3.7 has been the starting point of a very active line of research, composed of several hundreds papers – both theoretical and applied. In the next section, we will implicitly provide a general version of Theorem 3.6 (with the 1-Wasserstein distance replacing the total variation distance), whose proof relies only on the spectral properties of the Ornstein-Uhlenbeck generator $L$ and on the so-called $\Gamma$ calculus (see e.g. [BGL14]).

4 The Markov triple approach

In this section, we introduce a general framework for studying and generalizing the fourth moment phenomenon appearing in the statement of Corollary 3.7. The forthcoming approach was first introduced in [Led12] by M. Ledoux, and then further developed and generalizes in [ACP14, AMMP16].

4.1 Diffusive fourth moment structures

We start with definition of our general setup.

Definition 4.1. A diffusive fourth moment structure is a triple $(E, \mu, L)$ such that:

(a) $(E, \mu)$ is a probability space;

(b) $L$ is a symmetric unbounded operator defined on some dense subset of $L^2(E, \mu)$, that we denote by $\mathcal{D}(L)$ (the set $\mathcal{D}(L)$ is called the domain of $L$);

(c) the associated carré-du-champ operator $\Gamma$ is a symmetric bilinear operator, and is defined by

$$2\Gamma [X, Y] := L[XY] - XL[Y] - YL[X].$$

(4.1)
(d) the operator $L$ is **diffusive**, meaning that, for any $C^2$ function $\varphi: \mathbb{R} \to \mathbb{R}$, any $X \in D(L)$, it holds that $\varphi(X) \in D(L)$ and

$$L[\varphi(X)] = \varphi'(X)L[X] + \varphi''(X)\Gamma[X, X];$$ \hspace{1cm} (4.2)

Note that $L[1] = 0$ (by taking $\varphi = 1 \in C^2$). The latter property is equivalent to say that operator $\Gamma$ satisfies in the chain rule:

$$\Gamma[\varphi(X), X] = \varphi'(X)\Gamma[X, X];$$

(e) the operator $-L$ diagonalizes the space $L^2(E, \mu)$ with $sp(-L) = \mathbb{N}$, meaning that

$$L^2(E, \mu) = \bigoplus_{i=0}^{\infty} \text{Ker}(L + i\text{Id});$$

(f) for any pair of eigenfunctions $(X, Y)$ of the operator $-L$ associated with the eigenvalues $(p_1, p_2)$,

$$XY \in \bigoplus_{i \leq p_1 + p_2} \text{Ker}(L + i\text{Id}).$$ \hspace{1cm} (4.3)

In this context, we usually write $\Gamma[X]$ instead of $\Gamma[X, X]$ and $E$ denotes the integration against probability measure $\mu$.

**Remark 4.2.** (1) Property (d) together with symmetric property of the operator $L$ equip us with a functional calculus through the following fundamental integration by parts formula: for any $X, Y \in D(L)$ and $\varphi \in C^2$,

$$E[\varphi'(X)\Gamma[X, Y]] = -E[\varphi(X)L[Y]] = -E[YL[\varphi(X)]];$$ \hspace{1cm} (4.4)

(2) The results in this section can be stated under the weaker assumption that $sp(-L) = \{0 = \lambda_0 < \lambda_1, \cdots, \lambda_k < \cdots\} \subset \mathbb{R}^+$ is discrete. However, to keep a transparent presentation, we restrict ourself to the assumption $sp(-L) = \mathbb{N}$. The reader is referred to [ACP14] for further details.

(3) We point out that, by a recursive argument, assumption (4.3) yields that for any $X \in \bigoplus_{i \leq mp} \text{Ker}(L + i\text{Id})$ and any polynomial $P$ of degree $m$, we have

$$P(X) \in \bigoplus_{i \leq mp} \text{Ker}(L + i\text{Id}).$$ \hspace{1cm} (4.5)

(4) The eigenspaces of a diffusive fourth moment structure are **hypercontractive** (see [B94] for details and sufficient conditions), that is, there exists a constant $C(M, k)$ such that for any $X \in \bigoplus_{i \leq M} \text{Ker}(L + i\text{Id})$:

$$E(X^{2k}) \leq C(M, k) E(X^k).$$ \hspace{1cm} (4.6)

(5) Property (f) in the previous definition roughly implies that eigenfunctions of $L$ in a diffusive fourth moment structure behave like orthogonal polynomial with respect to multiplication.

For further details on our setup, we refer the reader to [BGL14] as well as [ACP14, AMMP16]. The next example describes some diffusive fourth moment structures. The reader can consult [AMMP16, Section 2.2] for two classical methods for building further diffusive fourth moment structures starting from known ones.
Example 4.3. (a) Finite-Dimensional Gaussian Structures: Let $d \geq 1$ and denote by $\gamma_d$ the $d$-dimensional standard Gaussian measure on $\mathbb{R}^d$. It is well known (see for example [BGL14]), that $\gamma_d$ is the invariant measure of the Ornstein-Uhlenbeck generator, defined for any test function $\varphi$ by

$$L\varphi(x) = \Delta \varphi - \sum_{i=1}^d x_i \partial_i \varphi(x). \quad (4.7)$$

Its spectrum is given by $-\mathbb{N}_0$ and the eigenspaces are of the form

$$\text{Ker}(L + k\text{Id}) = \left\{ \sum_{i_1 + i_2 + \cdots + i_d = k} \alpha(i_1, \cdots, i_d) \prod_{j=1}^d H_{i_j}(x_j) \right\},$$

where $H_n$ denotes the Hermite polynomial of order $n$. Since, eigenfunctions of $L$ are multivariate polynomials so it is straightforward to see that assumption (f) is also verified.

(b) Wiener space and isonormal processes: Letting $d \to \infty$ in the setup of the previous item (a) one recovers the infinite dimensional generator of the Ornstein-Uhlenbeck semigroup for isonormal processes, as defined in Section 3.1. It is easily verified in particular, by using (4.5) that $(\Omega, \mathcal{F}, L)$ is also a diffusive fourth moment structure.

(c) Laguerre Structure: Let $\nu \geq -1$, and $\pi_{1,\nu}(dx) = x^{\nu-1} e^{-x} \frac{1}{\Gamma(\nu)} 1_{(0,\infty)}(x) dx$ be the Gamma distribution with parameter $\nu$ on $\mathbb{R}_+$. The associated Laguerre generator is defined for any test function $\varphi$ (in dimension one) by:

$$L_{1,\nu}(\varphi) = x \varphi''(x) + (\nu + 1 - x) \varphi'(x). \quad (4.8)$$

By a classical tensorization procedure, we obtain the Laguerre generator in dimension $d$ associated with the measure $\pi_{d,\nu}(dx) = \pi_{1,\nu}(dx_1) \pi_{1,\nu}(dx_2) \cdots \pi_{1,\nu}(dx_d)$, where $x = (x_1, x_2, \cdots, x_d)$.

$$L_{d,\nu}(\varphi) = \sum_{i=1}^d \left( x_i \partial_i \varphi + (\nu + 1 - x_i) \partial_i \varphi \right) \quad (4.9)$$

It is also classical that (see for example [BGL14]) the spectrum of $L_{d,\nu}$ is given by $-\mathbb{N}_0$ and moreover that

$$\text{Ker}(L_{d,\nu} + k\text{Id}) = \left\{ \sum_{i_1 + i_2 + \cdots + i_d = k} \alpha(i_1, \cdots, i_d) \prod_{j=1}^d L_{i_j}(\nu)(x_j) \right\}, \quad (4.10)$$

where $L_{i_j}(\nu)$ stands for the Laguerre polynomial of order $n$ with parameter $\nu$ which is defined by

$$L_{n}(\nu)(x) = \frac{x^{-\nu} e^x}{n!} \frac{d^n}{dx^n} \left( e^{-x} x^{n+\nu} \right).$$

In the next subsection, we demonstrate how a diffusive fourth moment structure can be combined with the tools of $\Gamma$ calculus, in order to deduce substantial generalizations of Theorem 5.6.

4.2 Connection with $\Gamma$ Calculus

Throughout this section, we assume that $(E, \mu, L)$ is a diffusive fourth moment structure. Our principal aim is to prove an analogous fourth moment criterion to that of (3.12) for eigenfunctions of the operator $L$. To do this, we assume that $X \in \text{Ker}(L + q\text{Id})$ for some $q \geq 1$ with $E[X^2] = 1$. The arguments implemented in the proof will clearly demonstrate that requirements (d) and (f) in Definition 4.1 are the most crucial elements in order to establish our estimates.
Proposition 4.4. Let \( q \geq 1 \). Assume that \( X \in \text{Ker}(L + q\text{Id}) \) with \( \mathbb{E}[X^2] = 1 \). Then,

\[
\text{Var}(\Gamma[X]) \leq \frac{q^2}{3} \{ \mathbb{E}[X^4] - 3 \}.
\]

**Proof.** First note that by using integration by parts formula (4.4), we have \( \mathbb{E}[\Gamma[X]] = -\mathbb{E}[X\mathbb{L}X] = q\mathbb{E}[X^2] = q \). Secondly, by using the definition of the carré-du-champ operator \( \Gamma \) and the fact that \( \mathbb{L}X = -qX \), one easily verifies that

\[
\Gamma[X] - q = \frac{1}{2} (\mathbb{L} + 2q\text{Id}) (X^2 - 1).
\]

Next, taking into account properties (f) and (g) we can conclude that

\[
X^2 - 1 \in \bigoplus_{1 \leq i \leq 2q} \text{Ker}(\mathbb{L} + i\text{Id}).
\]

For the rest of the proof, we use the notation \( J_i \) to denote the projection of a square-integrable element \( X \) over the eigenspace \( \text{Ker}(\mathbb{L} + i\text{Id}) \). Now,

\[
\text{Var}(\Gamma[X]) = \mathbb{E}[(\Gamma[X] - q)^2] = \frac{1}{4} \mathbb{E}[(\mathbb{L} + 2q\text{Id})(X^2 - 1) \times (\mathbb{L} + 2q\text{Id})(X^2 - 1)]
\]

\[
= \frac{1}{4} \mathbb{E}[(\mathbb{L}(X^2 - 1))(\mathbb{L} + 2q\text{Id})(X^2 - 1)] + \frac{q}{2} \mathbb{E}[(X^2 - 1)(\mathbb{L} + 2q\text{Id})(X^2 - 1)]
\]

\[
= \frac{1}{4} \sum_{1 \leq i \leq 2q} (-i)(2q - i)\mathbb{E}[(J_i(X^2 - 1))^2] + \frac{q}{2} \mathbb{E}[(X^2 - 1)(\mathbb{L} + 2q\text{Id})(X^2 - 1)]
\]

\[
\leq \frac{q}{2} \mathbb{E}[(X^2 - 1)(\mathbb{L} + 2q\text{Id})(X^2 - 1)]
\]

\[
= q \mathbb{E}[(X^2 - 1)(\Gamma[X] - q)] = q \mathbb{E}[(X^2 - 1)\Gamma[X]]
\]

\[
= q \mathbb{E}[(\Gamma[\frac{X^3}{3} - X, X])] = -q \mathbb{E}[(\frac{X^3}{3} - X)\mathbb{L}X]
\]

\[
= q^2 \mathbb{E}[X(\frac{X^3}{3} - X)] = q^2 \mathbb{E}[\frac{X^4}{3} - X^2]
\]

\[
= \frac{q^2}{3} \{ \mathbb{E}[X^4] - 3 \},
\]

thus yielding the desired conclusion. \( \Box \)

In order to avoid some technicalities, we now present a quantitative bound in the 1-Wasserstein distance \( d_1 \) (and not in the more challenging total variation distance \( d_{TV} \)) for eigenfunctions of the operator \( \mathbb{L} \). This requires to adapt the Stein’s method machinery presented in Section 2 to our setting, as a direct application of the integration by part formula (4.4). The arguments below are borrowed in particular from [Led12, Proposition 1].

**Proposition 4.5.** Let \((E, \mu, \mathbb{L})\) be a diffusive fourth moment structure. Assume that \( X \in \text{Ker}(\mathbb{L} + q\text{Id}) \) for some \( q \geq 1 \) with \( \mathbb{E}[X^2] = 1 \). Let \( N \sim \mathcal{N}(0, 1) \). Then,

\[
d_1(X, N) \leq \frac{2}{q} \text{Var}(\Gamma[X])^{\frac{1}{2}}.
\]

**Proof.** For every function \( f \) of class \( C^1 \) on \( \mathbb{R} \), with \( ||f'||_\infty \leq 1 \) and \( f' \in \text{Lip}(2) \) according to Part (b) in Lemma 2.1, it is enough to show that

\[
\mathbb{E}[(f'(X) - Xf(X))] \leq \frac{2}{q} \text{Var}(\Gamma[X])^{\frac{1}{2}}.
\]
Since \( L X = -qX \), and diffusivity of the operator \( \Gamma \) together with integration by parts formula one can write that
\[
E \left[ f'(X) - X f(X) \right] = E \left[ f'(X) + \frac{1}{q} L(X) f(X) \right] = E \left[ f'(X) - \frac{1}{q} \Gamma[f(X), X] \right] = E \left[ f'(X) - \frac{1}{q} f'(X) \Gamma[X] \right] = \frac{1}{q} E \left[ f'(X) (q - \Gamma[X]) \right].
\]

Now, the claim follows at once by using the Cauchy-Schwarz inequality and noting that \( E[\Gamma[X]] = q E[X^2] = q \).

We end this section with the following general version of the fourth moment theorem for eigenfunctions of the operator \( L \), obtained by combining Propositions \[13\] and \[15\].

**Theorem 4.6.** Let \( (E, \mu, L) \) be a diffusive fourth moment structure. Assume that \( X \in \text{Ker}(L + q \text{Id}) \) for some \( q \geq 1 \) with \( E[X^2] = 1 \). Let \( N \sim \mathcal{N}(0, 1) \). Then,
\[
d_W(X, N) \leq \frac{2}{\sqrt{3}} \sqrt{E[X^4] - 3}.
\]

It follows that, if \( \{X_n\}_{n \geq 1} \) is a sequence of eigenfunctions in a fixed eigenspace \( \text{Ker}(L + q \text{Id}) \) where \( q \geq 1 \) and \( E[X_n^2] = 1 \) for all \( n \geq 1 \), then the following implication holds: \( E[X_n^4] \to 3 \) if and only if \( X_n \) converges in distribution towards the standard Gaussian random variable \( N \).

**Remark 4.7.** The fact that the condition \( E[X_n^4] \to 3 \) is necessary for convergence to Gaussian is a direct consequence of the hypercontractive estimate \((1.10)\).

### 4.3 Transport distance, Stein discrepancy and \( \Gamma \) Calculus

The general setting of the Markov triple together with \( \Gamma \) calculus provide a suitable framework to study functional inequalities such as the classical logarithmic Sobolev inequality or the celebrated Talagrand quadratic transportation cost inequality. For simplicity, here we restrict ourself to the setting of Wiener structure and the Gaussian measure to be our reference measure. The reader may consult references [LNP15, LNP16] for a presentation of the general setting, and [NPS14, NPS15] for some previous references connecting fourth moment theorems and entropic estimates.

Let \( d \geq 1 \), and \( d\gamma(x) = (2\pi)^{-\frac{d}{2}} e^{-\|x\|^2} dx \) be the standard Gaussian measure on \( \mathbb{R}^d \). Assume that \( d\nu = h d\gamma \) is a probability measure on \( \mathbb{R}^d \) with a (smooth) density function \( h : \mathbb{R}^d \to \mathbb{R}_+ \) with respect to the Gaussian measure \( \gamma \). Inspired from Gaussian integration by parts formula we introduce first the crucial notion of a Stein kernel \( \tau_\nu \) associated with the probability measure \( \nu \) and, then, the concept of Stein discrepancy.

**Definition 4.8.** (a) A measurable matrix-valued map \( \tau_\nu \) on \( \mathbb{R}^d \) is called a Stein kernel for the centered probability measure \( \nu \) if for every smooth test function \( \phi : \mathbb{R}^d \to \mathbb{R} \),
\[
\int_{\mathbb{R}^d} x \cdot \nabla \phi d\nu = \int_{\mathbb{R}^d} \langle \tau_\nu, \text{Hess}(\phi) \rangle_{\text{HS}} d\nu,
\]
where \( \text{Hess}(\phi) \) stands for the Hessian of \( \phi \), and \( \langle , \rangle_{\text{HS}} \) and \( \| , \|_{\text{HS}} \) denote the usual Hilbert-Schmidt scalar product and norm, respectively.

(b) The Stein discrepancy of \( \nu \) with respect to \( \gamma \) is defined as
\[
S(\nu, \gamma) = \inf \left( \int_{\mathbb{R}^d} \| \tau_\nu - \text{Id} \|_{\text{HS}}^2 d\nu \right)^{\frac{1}{2}}
\]
where the infimum is taken over all Stein kernels of \( \nu \), and takes the value \(+\infty\) if a Stein kernel for \( \nu \) does not exist.

We recall that the Stein kernel \( \tau_\nu \) is uniquely defined in dimension \( d = 1 \), and that unicity may fail in higher dimensions \( d \geq 2 \). Also, \( \tau_\nu = \text{Id}_{d \times d} \) the identity matrix. We further refer to [Fati18, CFP18] for existence of the Stein kernel in general settings. The interest of the Stein's discrepancy comes e.g. from the fact that – as a simple application of Stein’s method –

\[
d_{TV}(\nu, \gamma) \leq 2 \left( \int_{\mathbb{R}^d} |\tau_\nu - 1|^2 d\nu \right)^{\frac{1}{2}},
\]

yielding that \( d_{TV}(\nu, \gamma) \leq 2 S(\nu, \gamma) \); see [LNP15] for further details.

Next, we need the notion of Wasserstein distance. Let \( p \geq 1 \). Given two probability measure \( \nu \) and \( \mu \) on the Borel sets of \( \mathbb{R}^d \), whose marginals have finite moments of order \( p \), we define the \( p \)-Wasserstein distance between \( \nu \) and \( \mu \) as

\[
W_p(\nu, \mu) = \inf_{\pi} \left( \int_{\mathbb{R}^d \times \mathbb{R}^d} |x - y|^p d\pi(x, y) \right)^{\frac{1}{p}}
\]

where the infimum is taken over all probability measures \( \pi \) of \( \mathbb{R}^d \times \mathbb{R}^d \) with marginals \( \nu \) and \( \mu \); note that \( W_1 = d_W \), as defined in Section 2.

For a measure \( \nu = h\gamma \) with a smooth density function \( h \) on \( \mathbb{R}^d \), we recall that

\[
H(\nu, \gamma) := \int_{\mathbb{R}^d} h \log h d\gamma = \text{Ent}_\gamma(h)
\]

is the relative entropy of the measure \( \nu \) with respect to \( \gamma \), and

\[
I(\nu, \gamma) := \int_{\mathbb{R}^d} |\nabla h|^2 h d\gamma
\]

is the Fisher information of \( \nu \) with respect to \( \gamma \). After having established these notions, we can state two popular probabilistic/Entropic functional inequalities :

(i) [Logarithmic Sobolev inequality]: \( H(\nu, \gamma) \leq \frac{1}{2} I(\nu, \gamma) \).

(ii) [Talagrand quadratic transportation cost inequality]: \( W_2^2(\nu, \gamma) \leq 2 H(\nu, \gamma) \).

The next theorem is borrowed from [LNP15], and represents a significant improvement of the previous logarithmic Sobolev and Talagrand inequalities based on the use of Stein discrepancies: the techniques used in the proof are based on an interpolation argument along the Ornstein-Uhlenbeck semigroup. The reader is also referred to the recent works [Fati18, CFP18, Sam18] for related estimates of the Stein discrepancy based on the use of Poincaré inequalities, as well as on optimal transport techniques.

**Theorem 4.9.** Let \( dv = hd\gamma \) be a centered probability measure on \( \mathbb{R}^d \) with smooth density function \( h \) with respect to the standard Gaussian measure \( \gamma \).

(1) Then the following Gaussian HSI inequality holds:

\[
H(\nu, \gamma) \leq \frac{1}{2} S^2(\nu, \gamma) \log \left( 1 + \frac{I(\nu, \gamma)}{S^2(\nu, \gamma)} \right).
\]
(2) Assume further that $S(\nu, \gamma)$ and $H(\nu, \gamma)$ are both positive and finite. Then, the following Gaussian WSH inequality holds:

$$W_2(\nu, \gamma) \leq S(\nu, \gamma) \arccos\left( e^{-\frac{H(\nu, \gamma)}{2S(\nu, \gamma)}} \right).$$

The next section focuses on a discrete Markov structure for which exact fourth moment estimates are available.

5 Poisson measures and the fourth moment phenomenon

We will now describe a non-diffusive Markov triple for which a fourth moment result analogous to Proposition 4.5 holds. Such a Markov triple is associated with the space of square-integrable functionals of a Poisson measure on a general pair $(Z, \mathcal{Z})$, where $Z$ is a Polish space and $\mathcal{Z}$ is the associated Borel $\sigma$-field. The requirement that $Z$ is Polish – together with several other assumptions adopted in the present section – is made in order to simplify the discussion; the reader is referred to [DP18b, DVZ18] for statements and proofs in the most general setting. See also [L16, LP17] for an exhaustive presentation of tools of stochastic analysis for functionals of Poisson processes, as well as [PR16] for a discussion of the relevance of variational techniques in the framework of modern stochastic geometry.

5.1 Setup

Let $\mu$ be a non-atomic $\sigma$-finite measure on $(Z, \mathcal{Z})$, and set $\mathcal{Z}_\mu := \{ B \in \mathcal{Z} : \mu(B) < \infty \}$. In what follows, we will denote by

$$\eta = \{ \eta(B) : B \in \mathcal{Z} \}$$

a Poisson measure on $(Z, \mathcal{Z})$ with control (or intensity) $\mu$. This means that $\eta$ is a random field indexed by the elements of $\mathcal{Z}$, satisfying the following two properties: (i) for every finite collection $B_1, \ldots, B_m \in \mathcal{Z}$ of pairwise disjoint sets, the random variables $\eta(B_1), \ldots, \eta(B_m)$ are stochastically independent, and (ii) for every $B \in \mathcal{Z}$, the random variable $\eta(B)$ has the Poisson distribution with mean $\mu(B)$.$^2$ Whenever $B \in \mathcal{Z}_\mu$, we also write $\hat{\eta}(B) := \eta(B) - \mu(B)$ and denote by

$$\hat{\eta} = \{ \hat{\eta}(B) : B \in \mathcal{Z}_\mu \}$$

the compensated Poisson measure associated with $\eta$. Throughout this section, we assume that $\mathcal{F} = \sigma(\eta)$.

It is a well-known fact that one can regard the Poisson measure $\eta$ as a random element taking values in the space $N_\sigma = N_\sigma(Z)$ of all $\sigma$-finite point measures $\chi$ on $(Z, \mathcal{Z})$ that satisfy $\chi(B) \in \mathbb{N}_0 \cup \{ +\infty \}$ for all $B \in \mathcal{Z}$. Such a space is equipped with the smallest $\sigma$-field $\mathcal{N}_\sigma := \mathcal{N}_\sigma(Z)$ such that, for each $B \in \mathcal{Z}$, the mapping $N_\sigma \ni \chi \mapsto \chi(B) \in [0, +\infty]$ is measurable. In view of our assumptions on $Z$ and following e.g. [LP17] Section 6.1, throughout the paper we can assume without loss of generality that $\eta$ is proper, in the sense that $\eta$ can be $P$-a.s. represented in the form

$$\eta = \sum_{n=1}^{\eta(Z)} \delta_{X_n},$$

(5.1)

where $\{ X_n : n \geq 1 \}$ is a countable collection of random elements with values in $Z$ and where we write $\delta_z$ for the Dirac measure at $z$. Since we assume $\mu$ to be non-atomic, one has that $X_k \neq X_n$ for every $k \neq n$, $P$-a.s..
Now denote by $F(N,\sigma)$ the class of all measurable functions $f: N,\sigma \to \mathbb{R}$ and by $L^0(\Omega) := L^0(\Omega,\mathcal{F})$ the class of real-valued, measurable functions $F$ on $\Omega$. Note that, as $\mathcal{F} = \sigma(\eta)$, each $F \in L^0(\Omega)$ has the form $F = f(\eta)$ for some measurable function $f$. This $f$, called a representative of $F$, is $P_\eta$-a.s. uniquely defined, where $P_\eta = P \circ \eta^{-1}$ is the image measure of $P$ under $\eta$. Using a representative $f$ of $F$, one can introduce the add-one cost operator $D^+ = (D^+_z)_{z \in \mathbb{Z}}$ on $L^0(\Omega)$ as follows:

$$D^+_z F := f(\eta + \delta_z) - f(\eta), \quad z \in \mathbb{Z};$$

(5.2)

similarly, we define $D^-$ on $L^0(\Omega)$ as

$$D^- z F := f(\eta) - f(\eta - \delta_z), \quad \text{if } z \in \text{supp}(\eta), \quad \text{and } D^- z F := 0, \quad \text{otherwise},$$

(5.3)

where $\text{supp}(\chi) := \{ z \in \mathbb{Z} : \text{for all } A \in \mathcal{A} \text{ s.t. } z \in A: \chi(A) \geq 1 \}$ is the support of the measure $\chi \in N_\sigma$. We call $-D^-$ the remove-one cost operator associated with $\eta$. We stress that the definitions of $D^+ F$ and $D^- F$ are, respectively, $P \otimes \mu$-a.e. and $P$-a.s. independent of the choice of the representative $f$ — see e.g. the discussion contained in [DP18b, Section 2] and the references therein. Note that he operator $D^+$ can be straightforwardly iterated by setting $D^{(1)} := D^+$ and, for $n \geq 2$ and $z_1, \ldots, z_n \in \mathbb{Z}$ and $F \in L^0(\Omega)$, by recursively defining

$$D^{(n)}_{z_1,\ldots,z_n} F := D^+_{z_1} (D^{(n-1)}_{z_2,\ldots,z_n} F),$$

5.2 $L^1$ integration by parts

One of the most fundamental formulae in the theory of Poisson processes is the so-called Mecke formula stating that, for each measurable function $h : N,\sigma \times \mathbb{Z} \to [0, +\infty]$, the identity

$$\mathbb{E} \left[ \int \int h(\eta + \delta_z, z) \mu(dz) \right] = \mathbb{E} \left[ \int h(\eta, z) \eta(dz) \right]$$

(5.4)

holds true; see [LP17, Chapter 4] for a detailed discussion. Such a formula can be used in order to define an (approximate) integration by parts formula on the Poisson space.

For random variables $F, G \in L^0(\Omega)$ such that $D^+ F D^+ G \in L^1(P \otimes \mu)$, we define

$$\Gamma_0(F, G) := \frac{1}{2} \left\{ \int \left( D^+_z F D^+_z G \right) \mu(dz) + \int \left( D^-_z F D^-_z G \right) \eta(dz) \right\}$$

(5.5)

which verifies $\mathbb{E}[|\Gamma_0(F, G)|] < \infty$, and $\mathbb{E}[\Gamma_0(F, G)] = \mathbb{E}[\int \left( D^+_z F D^+_z G \right) \mu(dz)]$, in view of Mecke formula. The following statement, taken from [DP18b], can be regarded as an integration by parts formula in the framework of Poisson random measures, playing a role similar to that of Lemma 3.4 in the setting of Gaussian fields. It is an almost direct consequence of (5.4).

**Lemma 5.1** ($L^1$ integration by parts). Let $G, H \in L^0(\Omega)$ be such that

$$GD^+ H, \quad D^+ G D^+ H \in L^1(P \otimes \mu).$$

Then,

$$\mathbb{E} \left[ G \left( \int \left( D^+_z H \mu(dz) \right) - \int \left( D^-_z H \eta(dz) \right) \right) \right] = -\mathbb{E}[\Gamma_0(G, H)].$$

(5.6)

We will now focus on multiple Wiener-Itô integrals.
5.3 Multiple integrals

For an integer $p \geq 1$ we denote by $L^2(\mu^p)$ the Hilbert space of all square-integrable and real-valued functions on $\mathbb{Z}^p$ and we write $L^2_2(\mu^p)$ for the subspace of those functions in $L^2(\mu^p)$ which are $\mu^p$-a.e. symmetric. Moreover, for ease of notation, we denote by $\| \cdot \|_2$ and $\langle \cdot , \cdot \rangle_2$ the usual norm and scalar product on $L^2(\mu^p)$ for whatever value of $p$. We further define $L^2(\mu^p) := \mathbb{R}$. For $f \in L^2(\mu^p)$, we denote by $I_p(f)$ the multiple Wiener-Itô integral of $f$ with respect to $\eta$. If $p = 0$, then, by convention, $I_0(c) := c$ for each $c \in \mathbb{R}$. Now let $p, q \geq 0$ be integers. The following basic properties are proved e.g. in [L16], and are analogous to the properties of multiple integrals in a Gaussian framework, as discussed in Section 3.1:

1. $I_p(f) = I_p(\tilde{f})$, where $\tilde{f}$ denotes the canonical symmetrization of $f \in L^2(\mu^p)$;

2. $I_p(f) \in L^2(P)$, and $\mathbb{E}[I_p(f)I_q(g)] = \delta_{p,q} p! \langle \tilde{f}, \tilde{g} \rangle_2$, where $\delta_{p,q}$ denotes the Kronecker’s delta symbol.

As in the Gaussian framework of Section 3.1 for $p \geq 0$ the Hilbert space consisting of all random variables $I_p(f)$, $f \in L^2(\mu^p)$, is called the $p$-th Wiener chaos associated with $\eta$, and is customarily denoted by $C_p$. It is a crucial fact that every $F \in L^2(P)$ admits a unique representation

$$F = \mathbb{E}[F] + \sum_{p=1}^{\infty} I_p(f_p), \quad (5.7)$$

where $f_p \in L^2(\mu^p)$, $p \geq 1$, are suitable symmetric kernel functions, and the series converges in $L^2(P)$. Identity (5.7) is the analogous of relation (3.2), and is once again referred to as the chaotic decomposition of the functional $F \in L^2(P)$.

The multiple integrals discussed in this section also enjoy multiplicative properties similar to formula (5.5) above – see e.g. [L16, Proposition 5] for a precise statement. One consequence of such product formulae is that, if $F \in C_p$ and $G \in C_q$ are such that $FG$ is square-integrable, then

$$FG \in \bigoplus_{r=0}^{p+q} C_r, \quad (5.8)$$

which can be seen as a property analogous to (5.5).

5.4 Malliavin operators

We now briefly discuss Malliavin operators on the Poisson space.

1. The domain $\text{dom} D$ of the Malliavin derivative operator $D$ is the set of all $F \in L^2(P)$ such that the chaotic decomposition (5.7) of $F$ satisfies $\sum_{p=1}^{\infty} p p! \| f_p \|_2^2 < \infty$. For such an $F$, the random function $Z \ni z \mapsto D_z F \in L^2(P)$ is defined via

$$D_z F = \sum_{p=1}^{\infty} p I_{p-1}(f_p(z, \cdot)), \quad (5.9)$$

whenever $z$ is such that the series is converging in $L^2(P)$ (this happens a.e.-$\mu$), and set to zero otherwise; note that $f_p(z, \cdot)$ is an a.e. symmetric function on $\mathbb{Z}^{p-1}$. Hence, $DF = (D_z F)_{z \in Z}$ is indeed an element of $L^2(P \otimes \mu)$. It is well-known that $F \in \text{dom} D$ if and only if $D^+ F \in L^2(P \otimes \mu)$, and in this case

$$D_z F = D^+_z F, \quad P \otimes \mu - \text{a.e.}, \quad (5.10)$$
2. The domain \( \text{dom} \ L \) of the Ornstein-Uhlenbeck generator \( L \) is the set of those \( F \in L^2(P) \) whose chaotic decomposition \((5.7)\) verifies \( \sum_{p=1}^{\infty} p^2 p! \| f_p \|^2 < \infty \) (so that \( \text{dom} \ L \subset \text{dom} \ D \)) and, for \( F \in \text{dom} \ L \), one defines
\[
LF = - \sum_{p=1}^{\infty} p I_p(f_p). \tag{5.11}
\]
By definition, \( E[LF] = 0 \); also, from \((5.11)\) it is easy to see that \( L \) is symmetric, in the sense that
\[
E\left[ (LF)G \right] = E\left[ F(LG) \right]
\]
for all \( F, G \in \text{dom} \ L \). Note that, from \((5.11)\), it is immediate that the spectrum of \( -L \) is given by the nonnegative integers and that \( F \in \text{dom} \ L \) is an eigenfunction of \( -L \) with corresponding eigenvalue \( p \) if and only if \( F = I_p(f_p) \) for some \( f_p \in L^2(\mu^p) \), that is:
\[
C_p = \text{Ker}(L + pI).
\]
The following identity corresponds to formula (65) in \([L16]\): if \( F \in \text{dom} \ L \) is such that \( D^+ F \in L^1(P \otimes \mu) \), then
\[
LF = \int_{\mathbb{Z}} (D^+_z F) \mu(dz) - \int_{\mathbb{Z}} (D^-_z F) \eta(dz). \tag{5.12}
\]
3. For suitable random variables \( F, G \in \text{dom} \ L \) such that \( FG \in \text{dom} \ L \), we introduce the carré-du-champ operator \( \Gamma \) associated with \( L \) by
\[
\Gamma(F, G) := \frac{1}{2} (L(FG) - FLG - GLF). \tag{5.13}
\]
The symmetry of \( L \) implies immediately the crucial integration by parts formula
\[
E\left[ (LF)G \right] = E\left[ F(LG) \right] = -E\left[ \Gamma(F, G) \right]; \tag{5.14}
\]
we will see below that, for many random variables \( F, G \), relation \((5.14)\) is indeed the same as identity appearing in Lemma \((5.1)\).

The following result – proved in \([DP18b]\) – provides an explicit representation of the carré-du-champ operator \( \Gamma \) in terms of \( \Gamma_0 \), as introduced in \((5.5)\).

**Proposition 5.2.** For all \( F, G \in \text{dom} \ L \) such that \( FG \in \text{dom} \ L \) and
\[
DF, DG, FDG, GDF \in L^1(P \otimes \mu),
\]
we have that \( DF = D^+ F, DG = D^+ G, \) in such a way that \( DF DG = D^+ F D^+ G \in L^1(P \otimes \mu) \), and
\[
\Gamma(F, G) = \Gamma_0(F, G), \tag{5.15}
\]
where \( \Gamma_0 \) is defined in \((5.5)\).

One crucial consequence of such a result is that the operator \( \Gamma \) is not diffusive, in such a way that the triple \((\Omega, P, L)\) is not a diffusive fourth moment structure, such as the ones introduced in Definition \((4.1)\) it follows in particular that the machinery of Section \((4)\) cannot be directly applied.
5.5 Fourth moment theorems

Starting at least from the reference [PSTU10] (where Malliavin calculus and Stein’s method were first combined on the Poisson space), it has been an open problem for several years that of establishing a fourth moment bound similar to Theorem 4.6 on the Poisson space. As recalled above, the main difficulty in achieving such a result is the discrete nature of add-one and remove-one cost operators, preventing in particular the triple \((\Omega, P, L)\) from enjoying a diffusive property.

The next statement contains one of the main bounds proved in [DVZ18], and shows that a quantitative fourth moment bound is available on the Poisson space. Such a bound (which also has a multidimensional extension) is proved by a clever combination of Malliavin-type techniques with an infinitesimal version of the exchangeable pairs approach toward Stein’s method – see e.g. [CCS10].

**Theorem 5.3.** For \(q \geq 2\), let \(F = I_q(f_q)\) be a multiple integral of order \(q\) with respect to \(\hat{\eta}\), and assume that \(E[F^2] = 1\). Then,

\[
d_W(F, N) \leq \left( \sqrt{\frac{2}{\pi}} + \frac{4}{3} \right) \sqrt{E[F^4]} - 3.
\]

One should notice that the first bound of this type was proved in [DP18b] under slightly more restrictive assumptions; also, reference [DP18b] contains analogous bounds in the Kolmogorov distance, that are not achievable by using exchangeable pairs. In particular, one of the key estimates used in [DP18b] is the following remarkable equality and bound

\[
\frac{1}{2q} \int_Z E[|D^4 F|^2] \mu(dz) = \frac{3}{q} E[F^4] - E[F^4] \leq \frac{4q - 3}{2q} \left( E[F^4] - 3E[F^2]^2 \right),
\]

that are valid for every \(F \in C_q, q \geq 2\), such that the mapping \(z \mapsto D^4 F\) verifies some minimal integrability conditions. Further characterisations of fourth moment theorems on the Poisson space, in particular via the use of contraction operators and product formulae, are contained in [DP18c].

6 Malliavin-Stein for targets in the second Wiener chaos

In this section, we present a short overview on the recent development on Malliavin-Stein approach for target distributions in the second Gaussian Wiener chaos. We also formulate some important conjectures that will complement the approach. We adopt the same notation as in Section 3.1 above. Let \(W\) stands for an isonormal Gaussian process on a separable Hilbert space \(\mathcal{H}\). Recall that the elements in the second Wiener chaos are random variables having the general form \(F = I_2(f)\), with \(f \in \mathcal{H}^{\otimes 2}\). Notice that, if \(f = h \otimes h\), where \(h \in \mathcal{H}\) is such that \(\|h\|_\mathcal{H} = 1\), then using the multiplication formula one has \(I_2(f) \sim N^2 - 1\), where \(N \sim \mathcal{N}(0, 1)\). To any kernel \(f \in \mathcal{H}^{\otimes 2}\), we associate the following Hilbert-Schmidt operator

\[
A_f : \mathcal{H} \mapsto \mathcal{H} : g \mapsto f \otimes_1 g.
\]

We also write \(\{\alpha_{f,j}\}_{j \geq 1}\) and \(\{\epsilon_{f,j}\}_{j \geq 1}\), respectively, to indicate the (not necessarily distinct) eigenvalues of \(A_f\) and the corresponding eigenvectors. The next proposition gathers together some relevant properties of the elements of the second Wiener chaos associated with \(W\).

**Proposition 6.1** (See Section 2.7.4 in [NP12]). Let \(F = I_2(f), f \in \mathcal{H}^{\otimes 2}\), be a generic element of the second Wiener chaos of \(W\), and write \(\{\alpha_{f,k}\}_{k \geq 1}\) for the set of the eigenvalues of the associated Hilbert-Schmidt operator \(A_f\).

1. The following equality holds: \(F = \sum_{k \geq 1} \alpha_{f,k} (N_k^2 - 1)\), where \(\{N_k\}_{k \geq 1}\) is a sequence of i.i.d. \(\mathcal{N}(0, 1)\) random variables that are elements of the isonormal process \(W\), and the series converges in \(L^2\) and almost surely.
2. For any $r \geq 2$,

$$\kappa_r(F) = 2^{r-1}(r-1)! \sum_{k \geq 1} \alpha_{f,k}^r.$$

From now on, for simplicity, we consider the target distributions in the second Wiener chaos of the form

$$F_\infty = \sum_{i=1}^d \alpha_{\infty,i} (N_i^2 - 1)$$

(6.1)

where $N_i \sim \mathcal{N}(0, 1)$ are i.i.d, and the coefficients $\alpha_{\infty,i}$ are distinct. We also assume that $\mathbb{E}[F_\infty^2] = 1$. In the special case when $\alpha_{\infty,i} = 1$ for $1 \leq i \leq d$, the target random variable $F_\infty$ reduces to that of a centered chi-squared distribution with $d$ degree of freedom. The Malliavin-Stein approach has been successfully implemented in a series of papers [NP09a; DP18a; NP09b; NPR10]. The target random variables of the form (6.1) with $d = 2$, and $\alpha_{\infty,1} \times \alpha_{\infty,2} < 0$ belong to the so-called Variance–Gamma class of probability distributions. We refer to [Gau14; ET14] for development of Stein and Malliavin-Stein for the Variance–Gamma distributions. In this setting, the first obstacle for fully developing the Malliavin-Stein approach was the absence of a “suitable” Stein operator (meaning by that a differential operator with polynomial coefficients) for the candidate target distribution. This is the message of the next result. Also, the stability phenomenon of the weak convergence of the sequences in the second Wiener chaos is studied in [NP12] using tools in complex analysis.

**Theorem 6.2 (Stein characterization [AAPS17]).** Let $F_\infty$ belongs to the second Wiener chaos of the form (6.1). Consider polynomials $Q(x) = (P(x))^2 = \left(x \prod_{i=1}^d (x - \alpha_{\infty,i})\right)^2$ and, coefficients

$$a_l = \frac{P^{(l)}(0)}{l! 2^{l-1}}, \quad 1 \leq l \leq d+1,$$

$$b_l = \sum_{r=1}^{d+1} \frac{a_r}{(r-l+1)!} \kappa_{r-l+2}(F_\infty). \quad 2 \leq l \leq d+1.$$

Assume that $F$ is a general centered random variable living in a finite sum of Wiener chaoses (and hence smooth in the sense of Malliavin calculus). Then $F = F_\infty$ (equality in distribution) if and only if $\mathbb{E}[\mathcal{A}_\infty f(F)] = 0$ for all polynomials $f : \mathbb{R} \to \mathbb{R}$ where differential operator $\mathcal{A}_\infty$ of order $d$ is

$$\mathcal{A}_\infty f(x) := \sum_{l=2}^{d+1} \left(b_l - a_{l-1} x \right) f^{(d+2-l)}(x) - a_{d+1} x f(x).$$

(6.2)

The next essential conjecture formulates the non-Gaussian counterpart of the Stein’s Lemma 2.1. An affirmative answer will complete the Stein part of the approach in this delicate setting.

**Conjecture 6.3 (Stein Universality Lemma).** Let $\mathcal{H}$ denote an appropriate class of test functions. For every given test function $h \in \mathcal{H}$ consider the associated Stein equation

$$\mathcal{A}_\infty f(x) = h(x) - \mathbb{E} [h(F_\infty)].$$

(6.3)

Then equation (6.3) admits a bounded solution $f_h$ which is $d$ times differentiable and that $\|f_h^{(r)}\|_\infty < +\infty$ for all $r = 1, \ldots, d$ and the bounds are independent of the test function $h$.

The rest of the section is devoted to the first-ever quantitative estimates with target distributions in the second Wiener chaos. The first estimate is stated in terms of 2-Wasserstein transport distance $W_2$ (see Section 1.3 for definition). We highlight that the upper bound involves only finitely many cumulants, and therefore consistent with one of the ultimate goal of Malliavin-Stein approach. The second result is more general and rather intricate containing the iterated Gamma operators of the Malliavin calculus. See also [K17] for several related results of a quantitative nature.
Theorem 6.4 (AAPS17a). Let $F_n = \sum_{k=1}^{d\deg(Q)} \alpha_{n,k} (N_2 - 1)$ be a sequence belongs to the second Wiener chaos associated to the isonormal process $W$ so that $E[F_n^2] = 1$ for all $n \geq 1$. Assume that the target random variable $F_\infty$ as in (6.1) Define

$$\Delta(F_n) = \sum_{r=2}^{d\deg(Q)} \frac{Q_r(0)}{r!} \frac{\kappa_r(F_n)}{(r-1)!2^{r-1}}.$$ 

Then there exists a constant $C > 0$ depending only on target random variable $F_\infty$ (and independent of $n$) such that

$$W_2(F_n, F_\infty) \leq C \left( \sqrt{\Delta(F_n)} + \sum_{r=2}^{d+1} |\kappa_r(F_n) - \kappa_r(F_\infty)| \right) .$$

(6.4)

Example 6.5. Let $d = 2$ and $\alpha_{\infty,1} = -\alpha_{\infty,2} = 1/2$, then the target random variable $F_\infty$ $(= N_1 \times N_2$, where $N_1, N_2 \sim \mathcal{N}(0,1)$ are independent and equality holds in law) belongs to the class of Variance–Gamma distributions $VG_\kappa$ with parameters $\kappa = 3$ and $\theta = 0$. Then, [ET14 Corollary 5.10, part (a)] reads

$$d_W(F_n, F_\infty) \leq C \sqrt{\Delta(F_n) + 1/4 \kappa^2(F_n)} .$$

(6.5)

which is consistency with estimate (6.4). One has to note that for the target random variable $F_\infty$ it holds that $\kappa_3(F_\infty) = 0$. For a generalization of the estimate (6.5) to the higher moments and 2-Wasserstein distance see [AG17].

The next result provides a quantitative bound in the Kolmogorov distance. The proof relies on the classical Berry–Essen estimate in terms of bounding the difference of the characteristic functions. We recall that for two real-valued random variables $X$ and $Y$ the Kolmogorov distance is defined as

$$d_{Kol}(X,Y) := \sup_{x \in \mathbb{R}} |P(X \in (-\infty,x]) - P(Y \in (-\infty,x])| .$$

Theorem 6.6 (AMP16). Let $F_\infty$ be the target random variable in the second Wiener chaos of the form (6.1). Assume that $\{F_n\}_{n \geq 1}$ be a sequence of centered random elements living in a finite sum of the Wiener chaoses. Then there exists a constant $C$ (may depend on the sequence $F_n$ but independent of $n$) such that

$$d_{Kol}(F_n, F_\infty) \leq C \sqrt{\frac{\sum_{r=2}^{d+1} \left( |\kappa_r(F_n) - \kappa_r(F_\infty)| + \sum_{r=2}^{d+1} |\kappa_r(F_n) - \kappa_r(F_\infty)| \right)}{\text{Var} \left( \sum_{r=1}^{d+1} a_r \Gamma_{r-1}(F_n) \right) + \sum_{r=2}^{d+1} |\kappa_r(F_n) - \kappa_r(F_\infty)|}} .$$

(6.6)

Remark 6.7. We remark that when the sequence $\{F_n\}_{n \geq 1}$ appearing in Theorem 6.6 belongs to the second Wiener chaos, then [APP15] yields that

$$\text{Var} \left( \sum_{r=1}^{d+1} a_r \Gamma_{r-1}(F_n) \right) = \Delta(F_n)$$

where the quantity $\Delta(F_n)$ is as Theorem 6.4. As a result, the estimate (6.6) takes the form (compare with (6.3))

$$d_{Kol}(F_n, F_\infty) \leq C \sqrt{\Delta(F_n) + \sum_{r=2}^{d+1} |\kappa_r(F_n) - \kappa_r(F_\infty)|} .$$
We end the section with the following conjecture aiming to control the iterated Gamma operators of Malliavin calculus appearing in the RHS of the estimate (6.6) with finitely many cumulants. A successful path might go through first proving the estimate (6.8) where we name it as the $\Gamma_2$–Conjecture. Finally we point out that the estimate (6.8) has to be compared with the famous estimate $\text{Var}(\Gamma_1(F)) \leq C\kappa_4(F)$ in the normal approximation setting where $F$ is a chaotic random variable.

**Conjecture 6.8.** Let $F_\infty$ be the target random variable in the second Wiener chaos of the form (6.1). Assume that $F = I_q(f)$ be a chaotic random variable in the $q$th Wiener chaos with $q \geq 2$. Then there exists a general constant $C$ (may depend on $q$ and $d$) such that

$$\text{Var} \left( \sum_{r=1}^{d+1} a_r \Gamma_{r-1}(F) \right) \leq C\Delta(F). \quad (6.7)$$

In the particular case, when $d = 2$, and $\alpha_{\infty, 1} = -\alpha_{\infty, 2} = 1/2$, then the target random variable $F_\infty (= N_1 \times N_2$, where $N_1, N_2 \sim \mathcal{N}(0, 1)$ are independent, the estimate (6.7) boils down to that

$$\text{Var} (\Gamma_2(F) - F) \leq C \left\{ \frac{\kappa_6(F)}{3!} - 2\frac{\kappa_4(F)}{3!} + \kappa_2(F) \right\}. \quad (6.8)$$

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