

Partial functional quantization and generalized bridges

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In this article, we develop a new approach to functional quantization, which consists in discretizing only a finite subset of the Karhunen–Loève coordinates of a continuous Gaussian semimartingale X .

Using filtration enlargement techniques, we prove that the conditional distribution of X knowing its first Karhunen–Loève coordinates is a Gaussian semimartingale with respect to a larger filtration. This allows us to define the partial quantization of a solution of a stochastic differential equation with respect to X by simply plugging the partial functional quantization of X in the SDE.

Then we provide an upper bound of the L^p -partial quantization error for the solution of SDEs involving the $L^{p+\varepsilon}$ -partial quantization error for X , for $\varepsilon > 0$. The a.s. convergence is also investigated.

Incidentally, we show that the conditional distribution of a Gaussian semimartingale X , knowing that it stands in some given Voronoi cell of its functional quantization, is a (non-Gaussian) semimartingale. As a consequence, the functional stratification method developed in Corlay and Pagès [Functional quantization-based stratified sampling methods (2010) Preprint] amounted, in the case of solutions of SDEs, to using the Euler scheme of these SDEs in each Voronoi cell.

Keywords: Brownian bridge; Brownian motion; Cameron–Martin space; filtration enlargement; functional quantization; Gaussian process; Gaussian semimartingale; Karhunen–Loève; Ornstein–Uhlenbeck; stratification; vector quantization; Wiener integral

0. Introduction

Let $(\Omega, \mathcal{A}, \mathbb{P})$ be a probability space, and E a reflexive separable Banach space. The norm on E is denoted by $|\cdot|$. The quantization of a E -valued random variable X consists in its approximation by a random variable Y taking finitely many values. The resulting error of this discretization is measured by the L^p norm of $|X - Y|$. If we settle on a fixed maximum cardinal for $Y(\Omega)$, the minimization of the quantization error amounts to the minimization problem:

$$\min\{\| |X - Y| \|_p, Y: \Omega \rightarrow E \text{ measurable, } \text{card}(Y(\Omega)) \leq N\}. \quad (0.1)$$

A solution to (0.1) is an optimal quantizer of X . The corresponding quantization error is denoted by $\mathcal{E}_{N,p}(X) := \min\{\| |X - Y| \|_p, Y: \Omega \rightarrow E \text{ measurable, } \text{card}(Y(\Omega)) \leq N\}$. One usually drops the p subscript in the quadratic case ($p = 2$). This problem, initially investigated as a signal discretization method (Gersho and Gray [10]), has then been introduced in numerical probability to devise cubature methods (Pagès [23]) or to solve multidimensional stochastic control problems (Bally, Pagès and Printems [3]). Since the early 2000s, the infinite-dimensional setting has

been extensively investigated from both constructive numerical and theoretical viewpoints with a special attention paid to functional quantization, especially in the quadratic case (Luschgy and Pagès [18]) but also in some other Banach spaces (Wilbertz [28]). Stochastic processes are viewed as random variables taking values in functional spaces.

We now assume that X is a bi-measurable stochastic process on $[0, T]$ verifying $\int_0^T \mathbb{E}[|X_t|^2] dt < +\infty$, so that this can be viewed as a random variable valued in the separable Hilbert space $L^2([0, T])$. We assume that its covariance function Γ^X is continuous. In the seminal article on Gaussian functional quantization (Luschgy and Pagès [18]), it is shown that in the centered Gaussian case, linear subspaces U of $L^2([0, T])$ spanned by L^2 -optimal quantizers correspond to principal components of X . In other words, they are spanned by the first eigenvectors of the covariance operator of X . Thus, the quadratic optimal quantization of X involves its Karhunen–Loève eigensystem $(e_n^X, \lambda_n^X)_{n \geq 1}$. If Y is a quadratic N -optimal quantizer of X and $d^X(N)$ is the dimension of the subspace of $L^2([0, T])$ spanned by $Y(\Omega)$, the quadratic quantization error $\mathcal{E}_N^2(X)$ verifies

$$\mathcal{E}_N^2(X) = \sum_{j \geq m+1} \lambda_j^X + \mathcal{E}_N^2 \left(\bigotimes_{j=1}^m \mathcal{N}(0, \lambda_j^X) \right) \quad \text{for } m \geq d^X(N), \tag{0.2}$$

$$\mathcal{E}_N^2(X) < \sum_{j \geq m+1} \lambda_j^X + \mathcal{E}_N^2 \left(\bigotimes_{j=1}^m \mathcal{N}(0, \lambda_j^X) \right) \quad \text{for } 1 \leq m < d^X(N). \tag{0.3}$$

To perform optimal quantization, the decomposition is first truncated at a fixed order m and then the \mathbb{R}^m -valued Gaussian vector, constituted of the m first coordinates of the process on its Karhunen–Loève decomposition, is quantized. To reach optimality, we have to determine the optimal rank of truncation $d^X(N)$ (the quantization dimension) and the optimal $d^X(N)$ -dimensional quantizer corresponding to the first coordinates $\bigotimes_{j=1}^{d^X(N)} \mathcal{N}(0, \lambda_j^X)$. A sharply optimized database of quantizers of univariate and multivariate Gaussian distributions is available on the web site www.quantize.maths-fi.com (Pagès and Printems [25]) for download. Usual examples of such processes are the standard Brownian motion on $[0, T]$, the Brownian bridge on $[0, T]$, Ornstein–Uhlenbeck processes and the fractional Brownian motion. In Figure 1, we display the quadratic optimal N -quantizer of the fractional Brownian motion on $[0, 1]$ with Hurst exponent $H = 0.25$ and $N = 20$.

From a constructive viewpoint, the numerical computation of the optimal quantization or the optimal product quantization requires a numerical evaluation of the Karhunen–Loève eigenfunctions and eigenvalues, at least the very first terms. (As seen in Luschgy and Pagès [18,19], Luschgy, Pagès and Wilbertz [21], under rather general conditions on its eigenvalues, the quantization dimension of a Gaussian process increases asymptotically as the logarithm of the size of the quantizer. Hence, it is most likely that it is small. For instance, the quantization dimension of Brownian motion with $N = 10,000$ is 9.) The Karhunen–Loève decompositions of several usual Gaussian processes have a closed-form expression. This is the case for standard Brownian motion, Brownian bridge and Ornstein–Uhlenbeck processes. The case of Ornstein–Uhlenbeck processes is derived in Corlay and Pagès [6], in the general setting of an arbitrary initial vari-

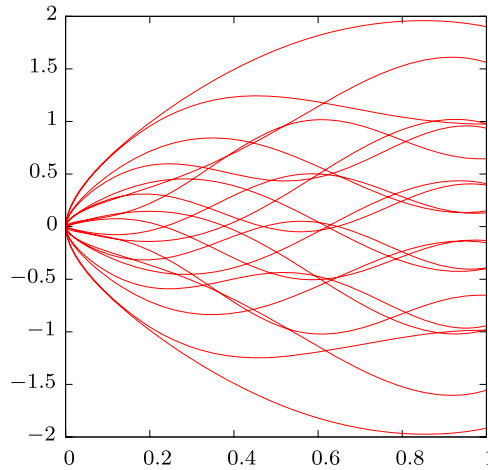


Figure 1. Quadratic N -optimal quantizer of the fractional Brownian motion on $[0, 1]$ with Hurst parameter $H = 0.25$ and $N = 20$. The quantization dimension is 3.

ance σ_0 . Another example of explicit Karhunen–Loève expansion is derived in Deheuvels and Martynov [8].

In the general case, no closed-form expression is available for the Karhunen–Loève expansion. For example, the K–L expansion of the fractional Brownian motion is not known. Yet, one can use numerical schemes to solve the corresponding eigenvalue problem. In Corlay [5], the so-called “Nyström method” is used to compute the first terms of the K–L decomposition of the fractional Brownian motion and to perform its optimal quantization.

In this article, we propose a new functional quantization scheme for a bi-measurable Gaussian process X , which consists in discretizing a finite subset of its Karhunen–Loève coordinates, instead of performing a full quantization. This *partial functional quantization* approach is motivated by two observations. The first one is that the conditional distribution of X knowing that it falls into a given L^2 Voronoi cell of its optimal quantizer is the crux of the recently developed functional stratification scheme (Corlay and Pagès [6]). It comes to conditioning the process with respect to its first Karhunen–Loève coordinates. This work provides a better justification of the functional stratification scheme of Corlay and Pagès [6]. The second observation is that one of the main purposes of the (full) functional quantization of X is to perform a quantization of the solution of a SDE with respect to X , when a stochastic integration with respect to X can be defined (see Pagès and Printems [24], Luschgy and Pagès [19], Pagès and Sellami [26]). As (full) functional quantizers of X will typically have bounded variations, one needs to add a correction term to the SDE. Eventually, this comes to plug the functional quantizer of X in the SDE written in the Stratonovich sense. In contrast, the partial quantization of X can be directly plugged into the SDE written in the Itô sense. We provide *a.s.* and L^p convergence results for this method.

The paper is organized as follows: Section 1 provides background on quantization-based cubature formulas which are needed for the following. In Section 2, we develop a notion of generalized bridge of a continuous Gaussian semimartingale which extends the generalized Brownian

bridge introduced in Alili [1]. We prove that under an additional hypothesis (\mathcal{H}), the generalized bridge of a continuous Gaussian semimartingale remains a Gaussian semimartingale with respect to a bigger filtration and we derive its canonical decomposition. (Let us mention the thorough study of the properties of Gaussian semimartingales available in Jain and Monrad [13].) A similar result is stated when conditioning by a Voronoi quantizer. We pay a particular attention to the special case of generalized bridges that we call Karhunen–Loève generalized bridges and which amounts to the conditioning of X by a finite subset of its K–L coordinates. Section 3 is devoted to the partial functional quantization of continuous Gaussian semimartingales and its application to the partial quantization of solutions of SDEs. We finally give L^p and a.s. convergence results for partially quantized SDEs.

1. Quantization-based cubature and related inequalities

The idea of quantization-based cubature method is to approach the probability distribution of the random variable X by the distribution of a quantizer Y of X . As Y is a discrete random variable, we can write $\mathbb{P}_Y = \sum_{i=1}^N p_i \delta_{y_i}$. If $F : E \rightarrow \mathbb{R}$ is a Borel functional,

$$\mathbb{E}[F(Y)] = \sum_{i=1}^N p_i F(y_i). \tag{1.1}$$

Hence, the weighted discrete distribution $(y_i, p_i)_{1 \leq i \leq N}$ of Y allows one to compute the sum (1.1). We review here some error bounds which can be derived when approaching $\mathbb{E}[F(X)]$ by (1.1). See Pagès and Printems [24] for detailed proofs.

1. If $X \in L^2$, Y a quantizer of X of size N and F is Lipschitz continuous, then

$$|\mathbb{E}[F(X)] - \mathbb{E}[F(Y)]| \leq [F]_{\text{Lip}} \|X - Y\|_2, \tag{1.2}$$

where $[F]_{\text{Lip}}$ is the Lipschitz constant of F . In particular, if $(Y_N)_{N \geq 1}$ is a sequence of quantizers such that $\lim_{N \rightarrow \infty} \|X - Y_N\|_2 = 0$, then the distribution $\sum_{i=1}^N p_i^N \delta_{x_i^N}$ of Y_N weakly converges to the distribution \mathbb{P}_X of X as $N \rightarrow \infty$.

This first error bound is a straightforward consequence of $|F(X) - F(Y)| \leq [F]_{\text{Lip}} |X - Y|$.

2. If Y is a stationary quantizer of X , that is, $Y = \mathbb{E}[X|Y]$, and F is differentiable with an α -Hölder differential DF for $\alpha \in (0, 1]$, that is, $|DF(u) - DF(v)|_{L(E)} \leq [DF]_\alpha |u - v|^\alpha$, for all $(u, v) \in E^2$ where $|\cdot|_{L(E)}$ is the operator norm on $L(E)$, then

$$|\mathbb{E}[F(X)] - \mathbb{E}[F(Y)]| \leq [DF]_\alpha \|X - Y\|_2^{1+\alpha}. \tag{1.3}$$

In the case where F has a Lipschitz continuous derivative ($\alpha = 1$), we have $[DF]_1 = [DF]_{\text{Lip}}$. For example, if F is twice differentiable and D^2F is bounded, then $[DF]_{\text{Lip}} = \|D^2F\|_\infty$.

This particular inequality comes from the Taylor expansion of F around X and the stationarity of Y .

3. If F is a semi-continuous¹ convex functional and Y is a stationary quantizer of X ,

$$\mathbb{E}[F(Y)] \leq \mathbb{E}[F(X)]. \tag{1.4}$$

This inequality is a straightforward consequence of the stationarity property and Jensen’s inequality.

$$\mathbb{E}[F(Y)] = \mathbb{E}[F(\mathbb{E}[X|Y])] \leq \mathbb{E}[\mathbb{E}[F(X)|Y]] = \mathbb{E}[F(X)].$$

2. Functional quantization and generalized bridges

2.1. Generalized bridges

Let $(X_t)_{t \in [0, T]}$ be a continuous centered Gaussian semimartingale starting from 0 on $(\Omega, \mathcal{A}, \mathbb{P})$ and \mathcal{F}^X its natural filtration. Fernique’s theorem ensures that $\int_0^T \mathbb{E}[X_t^2] dt < +\infty$ (see Janson [14]).

We aim here to compute the conditioning with respect to a finite family $\bar{Z}_T := (Z_T^i)_{i \in I}$ of Gaussian random variables, which are measurable with respect to $\sigma(X_t, t \in [0, T])$. ($I \subset \mathbb{N}$ is a finite subset of \mathbb{N}^* .) As in Alili [1] we settle on the case where $(Z_T^i)_{i \in I}$ are the terminal values of processes of the form $Z_t^i = \int_0^t f_i(s) dX_s, i \in I$, for some given finite set $\bar{f} = (f_i)_{i \in I}$ of $L^2_{\text{loc}}([0, T])$ functions. The *generalized bridge* for $(X_t)_{t \in [0, T]}$ corresponding to \bar{f} with end-point $\bar{z} = (z_i)_{i \in I}$ is the process $(X_t^{\bar{f}, \bar{z}})_{t \in [0, T]}$ that has the distribution

$$X^{\bar{f}, \bar{z}} \stackrel{\mathcal{L}}{\sim} \mathcal{L}(X | Z_T^i = z_i, i \in I). \tag{2.1}$$

For example, in the case where X is a standard Brownian motion with $|I| = 1, \bar{f} = \{f\}$ and $f \equiv 1$, this is the Brownian bridge on $[0, T]$. If X is an Ornstein–Uhlenbeck process, this is an Ornstein–Uhlenbeck bridge.

Let H be the Gaussian Hilbert space spanned by $(X_s)_{s \in [0, T]}$ and $H_{\bar{Z}_T}$ the closed subspace of H spanned by $(Z_T^i)_{i \in I}$. We denote by $H_{\bar{Z}_T}^\perp$ its orthogonal complement in H . Any Gaussian random variable G of H can be orthogonally decomposed into $G = \text{Proj}_{\bar{Z}_T}(G) \overset{\perp}{+} \text{Proj}_{\bar{Z}_T}^\perp(G)$, where $\text{Proj}_{\bar{Z}_T}$ and $\text{Proj}_{\bar{Z}_T}^\perp$ are the orthogonal projections on $H_{\bar{Z}_T}$ and $H_{\bar{Z}_T}^\perp$. ($\text{Proj}_{\bar{Z}_T}^\perp = Id_H - \text{Proj}_{\bar{Z}_T}$.) With these notation, $\mathbb{E}[G | (Z_T^i)_{i \in I}] = \text{Proj}_{\bar{Z}_T}(G)$.

Other definitions of generalized bridges exist in the literature, see, for example, Mansuy and Yor [22].

¹In the infinite-dimensional case, convexity does not imply continuity. In infinite-dimensional Banach spaces, a semi-continuity hypothesis is required for Jensen’s inequality. See Zapala [31] for more details.

2.2. The case of the Karhunen–Loève basis

As X is a continuous Gaussian process, it has a continuous covariance function (see Janson [14], Section VIII.3). We denote by $(e_i^X, \lambda_i^X)_{i \geq 1}$ its Karhunen–Loève eigensystem. Thus, if we define function f_i^X as the antiderivative of $-e_i^X$ that vanishes at $t = T$, that is, $f_i^X(t) = \int_t^T e_i^X(s) ds$, an integration by parts yields

$$\int_0^T X_s e_i^X(s) ds = \int_0^T f_i^X(s) dX_s. \tag{2.2}$$

In other words, with the notation of Section 2.1, we have $Z_T^i = \int_0^T X_s e_i^X(s) ds =: Y_i$, the i th Karhunen–Loève coordinate of X .

For some finite subset $I \subset \mathbb{N}^*$, we denote by $X^{I, \bar{y}}$ and call K – L generalized bridge the generalized bridge associated with functions $(f_i^X)_{i \in I}$ and with end-point $\bar{y} = (y_i)_{i \in I}$. This process has the distribution $\mathcal{L}(X|Y_i = y_i, i \in I)$.

In this case, the Karhunen–Loève expansion gives the decomposition

$$X = \underbrace{\sum_{i \in I} Y_i e_i^X}_{= \text{Proj}_{\bar{Z}_T}(X)} + \underbrace{\sum_{i \in \mathbb{N}^* \setminus I} \sqrt{\lambda_i^X} \xi_i e_i^X}_{= \text{Proj}_{\bar{Z}_T}^\perp(X)}, \tag{2.3}$$

where $(\xi_i)_{i \in \mathbb{N}^* \setminus I}$ are independent standard Gaussian random variables. This gives us the projections $\text{Proj}_{\bar{Z}_T}$ and $\text{Proj}_{\bar{Z}_T}^\perp$ defined in Section 2.1. It follows from (2.3) that a K–L generalized bridge is centered on $\mathbb{E}[X|Y_i = y_i, i \in I]$ and has the covariance function

$$\Gamma^{X|Y}(s, t) = \text{cov}(X_s, X_t) - \sum_{i \in I} \lambda_i^X e_i^X(s) e_i^X(t). \tag{2.4}$$

We have $\int_0^T \Gamma^{X|Y}(t, t) dt = \sum_{i \in \mathbb{N}^* \setminus I} \lambda_i^X$.

Moreover, thanks to decomposition (2.3), if $X^{I, \bar{y}}$ is a K–L generalized bridge associated with X with terminal values $\bar{y} = (y_i)_{i \in I}$, it has the same probability distribution as the process

$$\sum_{i \in I} y_i e_i^X(t) + X_t - \sum_{i \in I} \left(\int_0^T X_s e_i^X(s) ds \right) e_i^X(t).$$

This process is then the sum of a semimartingale and a non-adapted finite-variation process.

Let us stress the fact that the second term in the left-hand side of (2.3) is the corresponding K–L generalized bridge with end-point 0, that is, $\text{Proj}_{\bar{Z}_T}^\perp = X^{I, \bar{0}}$.

In Corlay and Pagès [6], an algorithm is proposed to exactly simulate marginals of a K–L generalized bridge with a linear additional cost to a prior simulation of $(X_{t_0}, \dots, X_{t_n})$, for some subdivision $0 = t_0 \leq t_1 \leq \dots \leq t_n = T$ of $[0, T]$. This was used for variance reduction issues. Note that the algorithm is easily extended to the case of (non-K–L) generalized bridges.

2.3. Generalized bridges as semimartingales

For a random variable L , we denote by $\mathbb{P}[\cdot|L]$ the conditional probability knowing L . We keep the notation and assumptions of previous sections. (X is a continuous Gaussian semimartingale starting from 0.) We consider a finite set $I \subset \{1, 2, \dots\}$ and $(f_i)_{i \in I}$ a set of bounded measurable functions. Let $X^{\bar{z}}$ be the generalized bridge associated with X with end-point $\bar{z} = (z_i)_{i \in I}$. For $i \in I$, $Z_t^i = \int_0^t f_i(s) dX_s$ and $\bar{Z}_t = (Z_t^i)_{i \in I}$.

Jirina’s theorem ensures the existence of a transition kernel

$$\nu_{\bar{Z}_T | ((X_t)_{t \in [0,s]})} : \mathcal{B}(\mathbb{R}^I) \times C^0([0, s], \mathbb{R}) \rightarrow \mathbb{R}_+,$$

corresponding to the conditional distribution $\mathcal{L}(\bar{Z}_t | (X_u)_{u \in [0,s]})$.

We now make the additional assumption (\mathcal{H}) that, for every $s \in [0, T)$ and for every $(x_u)_{u \in [0,s]} \in C^0([0, s], \mathbb{R})$, the probability measure $\nu_{\bar{Z}_T | (X_t)_{t \in [0,s]}}(d\bar{z}, (x_u)_{u \in [0,s]})$ is absolutely continuous with respect to the Lebesgue measure. We denote by $\Pi_{(x_u)_{u \in [0,s]}, T}$ its density. The covariance matrix of this Gaussian distribution on \mathbb{R}^I writes

$$Q(s, T) := \mathbb{E}\left[\left(\bar{Z}_T - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}\right]\right)\left(\bar{Z}_T - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}\right]^* \middle| (X_u)_{u \in [0,s]}\right].$$

If X is a martingale, we have $Q(s, T) = ((\int_s^T f_i(u) f_j(u) d\langle X \rangle_u))_{(i,j) \in I^2}$. We recall that a continuous centered semimartingale X is Gaussian if and only if $\langle X \rangle$ is deterministic (see, e.g., Revuz and Yor [27]). Hence, this additional hypothesis is equivalent to assume that

$$Q(s, T) \quad \text{is invertible for every } s \in [0, T). \tag{\mathcal{H}}$$

The following theorem follows from the same approach as the homologous result in Alili [1] for the Brownian case. It is extended to the case of a continuous centered Gaussian semimartingale starting from 0.

Theorem 2.1. *Under the (\mathcal{H}) hypothesis, for any $s \in [0, T)$, and for $\mathbb{P}_{\bar{Z}_T}$ -almost every $\bar{z} \in \mathbb{R}^I$, $\mathbb{P}[\cdot | \bar{Z}_T = \bar{z}]$ is equivalent to \mathbb{P} on \mathcal{F}_s^X and its Radon–Nikodym density is given by*

$$\frac{d\mathbb{P}[\cdot | \bar{Z}_T = \bar{z}]}{d\mathbb{P}} \Big|_{\mathcal{F}_s^X} = \frac{\Pi_{(X_u)_{u \in [0,s]}, T}(\bar{z})}{\Pi_{0,T}(\bar{z})}.$$

Proof. Consider F a real bounded \mathcal{F}_s^X -measurable random variable and $\phi : \mathbb{R}^I \rightarrow \mathbb{R}$ a bounded Borel function.

- On the one hand, preconditioning by \bar{Z}_T yields

$$\mathbb{E}[F \phi(\bar{Z}_T)] = \mathbb{E}[\mathbb{E}[F | \bar{Z}_T] \phi(\bar{Z}_T)] = \int_{\mathbb{R}^I} \phi(\bar{z}) \mathbb{E}[F | \bar{Z}_T = \bar{z}] \Pi_{0,T}(\bar{z}) d\bar{z}. \tag{2.5}$$

- On the other hand, as F is measurable with respect to \mathcal{F}_s^X , preconditioning with respect to \mathcal{F}_s^X yields

$$\mathbb{E}[F\phi(\bar{Z}_T)] = \mathbb{E}[F\mathbb{E}[\phi(\bar{Z}_T)|\mathcal{F}_s^X]] = \mathbb{E}\left[F \int_{\mathbb{R}^I} \phi(\bar{z}) \Pi_{(X_t)_{t \in [0,s]}, T}(\bar{z}) \, d\bar{z}\right].$$

Now, thanks to Fubini’s theorem

$$\mathbb{E}[F\phi(\bar{Z}_T)] = \int_{\mathbb{R}^I} \phi(\bar{z}) \mathbb{E}[F \Pi_{(X_t)_{t \in [0,s]}, T}(\bar{z})] \, d\bar{z}. \tag{2.6}$$

Identifying equations (2.5) and (2.6), we see that for $\mathbb{P}_{\bar{Z}_T}$ -almost surely $\bar{z} \in \mathbb{R}^I$ and for every real bounded \mathcal{F}_s^X -measurable random variable F ,

$$\mathbb{E}[F|\bar{Z}_T = \bar{z}] = \mathbb{E}\left[F \frac{\Pi_{(X_t)_{t \in [0,s]}, T}(\bar{z})}{\Pi_{0,T}(\bar{z})}\right]. \tag{2.7}$$

Equation (2.7) characterizes the Radon–Nikodym derivative of $\mathbb{P}[\cdot|\bar{Z}_T = \bar{z}]$ with respect to \mathbb{P} on \mathcal{F}_s^X . □

We now can use classical filtration enlargement techniques (Jacod [12], Jeulin [15], Yor [29]).

Proposition 2.1 (Generalized bridges as semimartingales). *Let us define the filtration $\mathcal{G}^{X,\bar{J}}$ by $\mathcal{G}_t^{X,\bar{J}} := \sigma(\bar{Z}_T, \mathcal{F}_t^X)$, the enlargement of the filtration \mathcal{F}^X corresponding to the above conditioning. We consider the stochastic process $D_s^{\bar{z}} := \frac{d\mathbb{P}[\cdot|\bar{Z}_T = \bar{z}]}{d\mathbb{P}}|_{\mathcal{F}_s^X} = \frac{\Pi_{(X_t)_{t \in [0,s]}, T}(\bar{z})}{\Pi_{0,T}(\bar{z})}$ for $s \in [0, T)$.*

Under the (\mathcal{H}) hypothesis, and the assumption that $D^{\bar{z}}$ is continuous, X is a continuous $\mathcal{G}^{X,\bar{J}}$ -semimartingale on $[0, T)$.

Proof. $D^{\bar{z}}$ is a strictly positive martingale on $[0, T)$ which is uniformly integrable on every interval $[0, t] \subset [0, T)$. Hence, as we assumed that it is continuous, we can write $D^{\bar{z}}$ as an exponential martingale $D_s^{\bar{z}} = \exp(L_s^{\bar{z}} - \frac{1}{2}\langle L^{\bar{z}} \rangle_s)$ with $L_t^{\bar{z}} = \int_0^t (D_s^{\bar{z}})^{-1} \, dD_s^{\bar{z}}$ (as $D_0^{\bar{z}} = 1$).

Now, as X is a continuous $(\mathcal{F}^X, \mathbb{P})$ -semimartingale, we write $X = V + M$ its canonical decomposition (under the filtration \mathcal{F}^X).

- Thanks to Girsanov theorem, $\tilde{M}^{\bar{z}} := M - \langle M, L^{\bar{z}} \rangle$ is a $(\mathcal{F}^X, \mathbb{P}[\cdot|\bar{Z}_T = \bar{z}])$ -martingale.
 - A consequence is that it is a $(\mathcal{G}^{X,\bar{J}}, \mathbb{P}[\cdot|\bar{Z}_T = \bar{z}])$ -martingale.
 - And thus $\tilde{M}^{\bar{z},T}$ is a $(\mathcal{G}^{X,\bar{J}}, \mathbb{P})$ -martingale.

For more preciseness on this, we refer to Ankirchner, Dereich and Imkeller [2], Theorem 3, where the proof is based on the notion of decoupling measure.

- Moreover, conditionally on \bar{Z}_T , V is still a finite-variation process V , and is adapted to $\mathcal{G}^{X,\bar{J}}$. □

Remark (Continuous modification). In Proposition 2.1, if one only assumes that $D^{\bar{z}}$ has a continuous modification $\mathcal{D}^{\bar{z}}$, then with each one of its continuous modifications is associated a continuous $\mathcal{G}^{X, \bar{f}}$ -semimartingale on $[0, T)$, and all these semimartingales are modifications of each other.

Proposition 2.2 (Continuity of $D^{\bar{z}}$). *If \mathcal{F}^X is a standard Brownian filtration, then $D^{\bar{z}}$ has a continuous modification.*

Proof. Consider $s \in [0, T)$. Under the (\mathcal{H}) hypothesis, the density $\Pi_{(X_u)_{u \in [0, s]}, T}$ writes

$$\begin{aligned} &\Pi_{(X_u)_{u \in [0, s]}, T}(\bar{z}) \\ &= (2\pi \det Q(s, T))^{-|I|/2} \\ &\quad \times \exp\left(\left(\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]\right) Q(s, T)^{-1} \left(\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]\right)^*\right). \end{aligned} \tag{2.8}$$

Let us define the stochastic process \bar{H} by $\bar{H}_s := \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]$. The so-defined process \bar{H} is a \mathcal{F}^X local martingale. Thanks the Brownian representation theorem, \bar{H} has a Brownian representation and has a continuous modification. The continuity of $s \mapsto \det Q(s, T)$ and $s \mapsto Q(s, T)^{-1}$ follows from the definition of $Q(s, T)$ and the continuity of \bar{H} (up to a modification). Hence, $D^{\bar{z}}$ has a continuous modification. \square

Remark.

- The measurability assumption with respect to a Brownian filtration is satisfied in the cases of Brownian bridge and Ornstein–Uhlenbeck processes.
- This hypothesis is not necessary so long as the continuity of the martingale $\bar{H}_s = \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]$ can be proved by any means.

2.3.1. *On the canonical decomposition*

With the same notation, and under the (\mathcal{H}) hypothesis, we can tackle the canonical decomposition of $X^{\bar{f}, \bar{z}}$. We have

$$L_t^{\bar{z}} = \int_0^t \frac{d\Pi_{(X_u)_{u \in [0, s]}, T}(\bar{z})}{\Pi_{(X_u)_{u \in [0, s]}, T}(\bar{z})}$$

and

$$\begin{aligned} &\ln(\Pi_{(X_u)_{u \in [0, s]}, T}(\bar{z})) \\ &= -\frac{|I|}{2} \ln(2\pi \det Q(s, T)) \\ &\quad - \frac{1}{2} \left(\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]\right) Q(s, T)^{-1} \left(\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, s]}]\right)^*. \end{aligned}$$

Using that for a positive continuous semimartingale S , $d \ln S = \frac{dS}{S} - \frac{1}{2} d\langle \frac{1}{S} \cdot S \rangle$, we obtain

$$\begin{aligned} & \frac{d\Pi_{(X_u)_{u \in [0,s]}, T}(\bar{z})}{\Pi_{(X_u)_{u \in [0,s]}, T}(\bar{z})} \\ &= d \ln(\Pi_{(X_u)_{u \in [0,s]}, T}(\bar{z})) + \left(\begin{array}{c} \text{finite-variation} \\ \text{process} \end{array} \right) \\ &= -\frac{1}{2} d\left((\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}]) Q(s, T)^{-1} (\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}])^* \right) + (\text{f.-v. p.}) \\ &= (d\mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}]) Q(s, T)^{-1} (\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}])^* + (\text{f.-v. p.}). \end{aligned}$$

Hence,

$$d\langle X, L^{\bar{z}} \rangle_s = d\langle X, \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, \cdot]}] \rangle_s Q(s, T)^{-1} (\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}])^*.$$

This expression can be further simplified in the two following cases:

- In the case where X is a martingale, owing to the definition of Z_j , we have $\forall j \in I$, $\mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}] = \int_0^s f_j(u) dX_u$ so that

$$\begin{aligned} d\langle X, L^{\bar{z}} \rangle_s &= (\bar{f}(s) Q(s, T)^{-1} (\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}])^*) d\langle X \rangle_s \\ &= \sum_{i \in I} f_i(s) \sum_{j \in I} (Q(s, T)^{-1})_{ij} (z_j - \mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}]) d\langle X \rangle_s. \end{aligned} \tag{2.9}$$

As a consequence, $M - \int_0^\cdot \sum_{i \in I} f_i(s) \sum_{j \in I} (Q(s, T)^{-1})_{ij} (z_j - \mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}]) d\langle X \rangle_s$ is a $(\mathcal{G}^{X, \bar{f}}, \mathbb{P}[\cdot | \bar{Z}_T = \bar{z}])$ -martingale. We have recovered Alili's result on the generalized Brownian bridge (Alili [1]).

- In the case where the Gaussian semimartingale X is a Markov process, for every $j \in I$ there exists $g_j \in L^2([0, T])$ such that $\mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}] = \int_0^s f_j(u) dX_u + g_j(s) X_s$. Indeed,

$$\mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}] = \int_0^s f_j(u) dX_u + \underbrace{\mathbb{E}\left[\int_s^T f_j(u) dX_u | (X_u)_{u \in [0,s]}\right]}_{=: g_j(s) X_s}.$$

Hence, if one assumes that $(g_j)_{j \in I}$ are finite-variation functions (which is the case when X is an Ornstein–Uhlenbeck process or a Brownian bridge), we have $d\langle X, \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0, \cdot]}] \rangle_s = (\bar{f}(s) + \bar{g}(s)) d\langle X \rangle_s$, and thus

$$\begin{aligned} d\langle X, L^{\bar{z}} \rangle_s &= ((\bar{f}(s) + \bar{g}(s)) Q(s, T)^{-1} (\bar{z} - \mathbb{E}[\bar{Z}_T | (X_u)_{u \in [0,s]}])^*) d\langle X \rangle_s \\ &= \sum_{i \in I} (f_i(s) + g_i(s)) \sum_{j \in I} (Q(s, T)^{-1})_{ij} (z_j - \mathbb{E}[Z_T^j | (X_u)_{u \in [0,s]}]) d\langle X \rangle_s. \end{aligned}$$

Example (Standard Brownian bridge). In the case where $X = W$ is a standard Brownian motion with $|I| = 1$, $\bar{f} = \{f\}$ and $f \equiv 1$, $Z_t = W_t$ and $W^{\bar{f}, \bar{z}}$ is a standard Brownian bridge. We have $Q(s, T)^{-1} = \frac{1}{T-s}$ and

$$d\langle X, L^z \rangle_t = \frac{1}{T-t} (z - \mathbb{E}[W_T | (W_u)_{u \in [0,t]}]) dt = \frac{z - W_t}{T-t} dt.$$

Thus,

$$dW_t = \frac{z - W_t}{T-t} dt + \underbrace{\frac{W_t}{T-t} dt + dW_t}_{(\mathcal{G}^{X, \bar{f}}, \mathbb{P}[\cdot | W_T = z])\text{-martingale}}.$$

The martingale part happens to be a $(\mathcal{G}^{X, \bar{f}}, \mathbb{P}[\cdot | W_T = z])$ -standard Brownian motion, thanks to Lévy’s characterization of the Brownian motion. Thus, we have retrieved the classical SDE of the Brownian bridge.

2.3.2. Generalized bridges and functional stratification

With the same notation, we set $\widehat{Z}^\Gamma = \text{Proj}_\Gamma(\bar{Z}_T) = \sum_{i=1}^N \gamma_i \mathbf{1}_{C_i}(\bar{Z}_T)$ a stationary quantizer of \bar{Z}_T (where $\Gamma = \{\gamma_1, \dots, \gamma_N\}$ and $C = \{C_1, \dots, C_N\}$ are, respectively, the associated knots and Voronoi partition).

Proposition 2.3 (Stratification). *Under the (\mathcal{H}) hypothesis, for any $s \in [0, T]$, for any $k \in \{1, \dots, N\}$, $\mathbb{P}[\widehat{Z}^\Gamma = \gamma_k] > 0$ and the conditional probability $\mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k]$ is equivalent to \mathbb{P} on \mathcal{F}_s^X .*

Proof. Obviously, if $A \in \mathcal{F}_s^X$ is such that $\mathbb{P}[A] = 0$, we have $\mathbb{P}[A | \widehat{Z}^\Gamma = \gamma_k] = 0$. Conversely, $B \in \mathcal{F}_s^X$ satisfies $\mathbb{P}[B | \widehat{Z}^\Gamma = \gamma_k] = 0$, then pre-conditioning by \bar{Z}_T , we get $\mathbb{E}[\mathbb{E}[\mathbf{1}_B | \bar{Z}_T] | \widehat{Z}^\Gamma = \gamma_k] = 0$. Thus, $\int_{\bar{z} \in C_k} \mathbb{P}[B | \bar{Z}_T = \bar{z}] d\mathbb{P}_{\bar{Z}_T}(\bar{z}) = 0$. Hence $\mathbb{P}[B | \bar{Z}_T = \bar{z}] = 0$ for $\mathbb{P}_{\bar{Z}_T}$ -almost every $\bar{z} \in C_k$.

Since $\mathbb{P}_{\bar{Z}_T}[C_k] > 0$, there exists at least one element $\bar{z} \in C_k$ such that $\mathbb{P}[B | \bar{Z}_T = \bar{z}] = 0$. Now thanks to Theorem 2.1, $\mathbb{P}[B] = 0$. □

Proposition 2.4 (Stratification). *Let us define the filtration $\mathcal{G}^{X, \Gamma}$ by $\mathcal{G}_t^{X, \Gamma} := \sigma(\mathcal{F}_t^X, \widehat{Z}^\Gamma)$, the enlargement of \mathcal{F}^X corresponding to the conditioning with respect to \widehat{Z}^Γ . For $k \in \{1, \dots, N\}$, we consider the stochastic process $D_s^{\gamma_k} := \frac{d\mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k]}{d\mathbb{P}} \Big|_{\mathcal{F}_s^X}$ for $s \in [0, T]$.*

Under the (\mathcal{H}) hypothesis, and the assumption that D^{γ_k} is continuous, the conditional distribution $\mathcal{L}(X | \widehat{Z}^\Gamma)$ of X knowing in which Voronoi cell \bar{Z}_T falls, is the probability distribution of a $\mathcal{G}^{X, \Gamma}$ -semimartingale on $[0, T]$.

Proof. Using that $\mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k]$ is equivalent to \mathbb{P} on \mathcal{F}_s^X , thanks to Proposition 2.3, we can *mutatis mutandis* use the same arguments as for Proposition 2.1, $\mathbb{P}[\cdot | \bar{Z}_T = \bar{z}]$ being replaced by $\mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k]$.

D^{γ_k} is a strictly positive martingale on $[0, T)$ uniformly integrable on every $[0, t] \subset [0, T)$. Hence, as D^{γ_k} is continuous by hypothesis, it is an exponential martingale $D_s^{\gamma_k} = \exp(L_s^{\gamma_k} - \frac{1}{2}\langle L^{\gamma_k} \rangle_s)$, with $L_t^{\gamma_k} = \int_0^t (D_s^{\gamma_k})^{-1} dD_s^{\gamma_k}$ (as $D_0^{\gamma_k} = 1$). Now, as X is a continuous $(\mathcal{F}^X, \mathbb{P})$ -semimartingale, we write $X = V + M$ its canonical decomposition (under the filtration \mathcal{F}^X).

- Thanks to Girsanov theorem, $\tilde{M}^{\gamma_k} := M - \langle M, L^{\gamma_k} \rangle$ is a $(\mathcal{F}^X, \mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k])$ -martingale. As a consequence, it is a $(\mathcal{G}^{X, \Gamma}, \mathbb{P}[\cdot | \widehat{Z}^\Gamma = \gamma_k])$ -martingale and thus $\tilde{M}^{\widehat{Z}^\Gamma}$ is a $(\mathcal{G}^{X, \Gamma}, \mathbb{P})$ -martingale.
- Moreover, conditionally to \widehat{Z}^Γ , V is still a finite-variation process V , and is adapted to $\mathcal{G}^{X, \Gamma}$. □

Proposition 2.5 (Continuity of D^{γ_k}). *If \mathcal{F}^X is a Brownian filtration, then D^{γ_k} has a continuous modification.*

Proof. By definition, D^{γ_k} is a \mathcal{F}^X -local martingale on $[0, T]$. The conclusion is a straightforward consequence of the Brownian representation theorem. □

Considering the partition of $L^2([0, T])$ corresponding to the Voronoi cells of a functional quantizer of X , the last two propositions show that the conditional distribution of the X in each Voronoi cell (strata) is a Gaussian semimartingale with respect to its own filtration. This allows us to define the corresponding functional stratification of the solutions of stochastic differential equations driven by X .

In Corlay and Pagès [6], an algorithm is proposed to simulate the conditional distribution of the marginals $(X_{t_0}, \dots, X_{t_n})$ of X for a given subdivision $0 = t_0 < t_1 < \dots < t_n = T$ of $[0, T]$ conditionally to a given Voronoi cell (strata) of a functional quantization of X . The simulation complexity has an additional linear complexity to an unconditioned simulation of $(X_{t_0}, \dots, X_{t_n})$. We refer to Corlay and Pagès [6] for more details.

To deal with the solution of a SDE, it was proposed in Corlay and Pagès [6] to simply plug these marginals in the Euler scheme of the SDE. Proposition 2.4 now shows that this amounts to simulate the Euler scheme of the SDE driven by the corresponding (non-Gaussian) semimartingale.

2.4. About the (\mathcal{H}) hypothesis

2.4.1. The martingale case

In the case where X is a continuous Gaussian martingale, the matrix $Q(s, t)$ defined in Section 2.3 writes $Q(s, t) = ((\int_s^t f_i(u) f_j(u) d\langle X \rangle_u))_{(i, j) \in I^2}$.

For $1 \leq s < t \leq T$, the map $(\cdot | \cdot) : (f, g) \mapsto \int_s^t f(u)g(u) d\langle X \rangle_u$ defines a scalar product on $L^2([s, t], d\langle X \rangle)$. Hence, $Q(s, t)$ is the Gram matrix of the vectors of $L^2([s, t], d\langle X \rangle)$ defined by the restrictions to $[s, t]$ of the functions $(f_i)_{i \in I}$. Thus, it is invertible if and only if these restrictions form a linearly independent family of $L^2([s, t], d\langle X \rangle)$. (Another consequence, is that if $Q(s, t)$ is invertible for some $0 \leq s < t \leq T$, then for every (u, v) such that $[s, t] \subset [u, v]$, $Q(u, v)$ is invertible.)

For instance, if X is a standard Brownian motion on $[0, T]$, the functions $(f_i^X)_{i \in I}$ (associated with the Karhunen–Loève decomposition) are trigonometric functions with strictly different frequencies. Hence, they form a linearly independent family of continuous functions on every non-empty interval $[s, T] \subset [0, T]$. Moreover, the measure $d\langle X \rangle$ is proportional to the Lebesgue measure on $[0, T]$ and thus $Q(s, T)$ is invertible for any $s \in [0, T]$. Hence, the (\mathcal{H}) hypothesis is fulfilled in the case of K–L generalized bridges of the standard Brownian motion.

2.4.2. Standard Brownian bridge and Ornstein–Uhlenbeck processes

Brownian bridge and the Ornstein–Uhlenbeck process are not martingales. Hence, this criterion is not sufficient and the invertibility of matrix $Q(s, T)$ has to be proved by other means.

Following from the definitions of $Q(s, T)$ and \bar{Z}_T , in the case of the K–L generalized bridge

$$\begin{aligned}
 Q(s, T)_{ij} &= \mathbb{E} \left[\left(\int_s^T f_i^X(u) dX_u - \mathbb{E} \left[\int_s^T f_i^X(u) dX_u \mid (X_u)_{u \in [0, s]} \right] \right) \right. \\
 &\quad \times \left. \left(\int_s^T f_j^X(u) dX_u - \mathbb{E} \left[\int_s^T f_j^X(u) dX_u \mid (X_u)_{u \in [0, s]} \right] \right)^* \mid (X_u)_{u \in [0, s]} \right] \quad (2.10) \\
 &= \text{cov} \left(\int_s^T f_i^X(u) dX_u^{(s)}, \int_s^T f_j^X(u) dX_u^{(s)} \right),
 \end{aligned}$$

where $(X_u^{(s)})_{u \in [s, T]}$ has the conditional distribution of X knowing $(X_u)_{u \in [0, s]}$.

- When X is a standard Brownian bridge on $[0, T]$, $(X_u^{(s)})_{u \in [s, T]}$ is a Brownian bridge on $[s, T]$, starting from X_s and arriving at 0. It is the sum of an affine function and a standard centered Brownian bridge on $[s, T]$.
- When X is a centered Ornstein–Uhlenbeck process, $(X_u^{(s)})_{u \in [s, T]}$ is an Ornstein–Uhlenbeck process on $[s, T]$ starting from X_s , with the same mean reversion parameter as X . It is also the sum of a deterministic function and an Ornstein–Uhlenbeck process starting from 0.

As a consequence, in these two cases, the quantity $\text{cov}(\int_s^T f_i^X(u) dX_u^{(s)}, \int_s^T f_j^X(u) dX_u^{(s)})$ can be computed by plugging either a centered Brownian bridge on $[s, T]$ or an Ornstein–Uhlenbeck starting from 0 instead of $X^{(s)}$ in equation (2.10). This means that $Q(s, T)$ is the Gram matrix of the random variables $(\int_s^T f_i^X(u) dG_u)_{i \in I}$, where the centered Gaussian process $(G_u)_{u \in [s, T]}$ is either a standard Brownian bridge on $[s, T]$ or an Ornstein–Uhlenbeck process starting from 0 at s . Thus, it is singular if and only if there exists $(\alpha_i)_{i \in I} \neq 0$ in \mathbb{R}^I such that

$$\int_s^T \underbrace{\left(\sum_{i \in I} \alpha_i f_i^X(u) \right)}_{=: g(u)} dG_u = 0 \quad \text{a.s.} \quad (2.11)$$

The case of Brownian bridge. In the case where X is the standard Brownian bridge on $[0, T]$, functions $(f_i^X)_{i \in I}$ are C^∞ functions and G is a standard Brownian bridge on $[s, T]$. An integration by parts gives $\int_s^T G_s g'(s) ds = 0$ a.s. and thus $g' \equiv 0$ on (s, T) and thus g is constant

on $[s, T]$. The functions $(f_i^X)_{i \in I}$ form a linearly independent set of functions and, as they are trigonometric functions with different frequencies, they clearly do not span constant functions, so that equation (2.11) yields $\alpha_1 = \dots = \alpha_n = 0$. Hence, the (\mathcal{H}) hypothesis is fulfilled in the case of K - L generalized bridges of the standard Brownian bridge.

The case of Ornstein–Uhlenbeck processes. In the case where X is an Ornstein–Uhlenbeck process on $[0, T]$, G is an Ornstein–Uhlenbeck process on $[s, T]$ starting from 0. The injectivity property of the Wiener integral related to the Ornstein–Uhlenbeck process stated in Proposition 2.6 below, applied on $[s, T]$, shows that equation (2.11) amounts to $g \stackrel{L^2([s, T], dt)}{=} 0$ and thus

$$\sum_{i \in I} \alpha_i f_i^X \stackrel{L^2([s, T], dt)}{=} 0. \tag{2.12}$$

Again, as $(f_i^X)_{i \in I}$ are linearly independent, we have $\alpha_1 = \dots = \alpha_n = 0$. Hence, the (\mathcal{H}) hypothesis is fulfilled in the case of K - L generalized bridges of the Ornstein–Uhlenbeck processes.

Proposition 2.6 (Injectivity of the Wiener integral related to centered Ornstein–Uhlenbeck processes). *Let G be an Ornstein–Uhlenbeck process defined on $[0, T]$ by the SDE*

$$dG_t = -\theta G_t dt + \sigma dW_t \quad \text{with } \sigma > 0 \text{ and } \theta > 0,$$

where W is a standard Brownian motion and $G_0 \stackrel{\mathcal{L}}{\sim} \mathcal{N}(0, \sigma_0^2)$ is independent of W .

If $g \in L^2([0, T])$, then we have

$$\int_0^T g(s) dG_s = 0 \quad \Leftrightarrow \quad g \stackrel{L^2([0, T])}{=} 0.$$

Proof. The solution of the Ornstein–Uhlenbeck SDE is

$$G_t = \underbrace{G_0 e^{-\theta t}}_{\text{independent of } W} \quad \perp \quad \underbrace{\int_0^t \sigma e^{\theta(s-t)} dW_s}_{=: G_t^0}.$$

Hence, we have

$$\int_0^T g(s) dG_s = -\theta G_0 \int_0^T g(s) e^{-\theta s} ds \quad \perp \quad \int_0^T g(s) dG_s^0.$$

Thus, by independence, if $\int_0^T g(s) dG_s = 0$ then $\int_0^T g(s) dG_s^0 = 0$. This means that we only have to prove the proposition in the case of an Ornstein–Uhlenbeck process starting from 0.

We now assume that $\sigma_0^2 = 0$ and we temporarily make the additional assumption that $\theta T < \frac{4}{3}$. If $g \in L^2([0, T])$ and $\int_0^T g(s) dG_s = 0$, then $\theta \int_0^T g(s) G_s ds = \sigma \int_0^T g(s) dW_s$, and thus, if Γ^{OU} denotes the covariance function of G ,

$$\theta^2 \int_0^T \int_0^T g(s) g(t) \Gamma^{\text{OU}}(s, t) ds dt = \sigma^2 \int_0^T g(s)^2 ds. \tag{2.13}$$

Applying Schwarz’s inequality twice, we get

$$\int_0^T \int_0^T g(s)g(t)\Gamma^{\text{OU}}(s, t) \, ds \, dt \leq \int_0^T g(s)^2 \, ds \sqrt{\int_0^T \int_0^T (\Gamma^{\text{OU}}(s, t))^2 \, ds \, dt}.$$

Hence, provided that

$$\int_0^T \int_0^T (\Gamma^{\text{OU}}(s, t))^2 \, ds \, dt < \frac{\sigma^4}{\theta^4}, \tag{2.14}$$

equality (2.13) implies $\int_0^T g(s)^2 \, ds = 0$.

Now, we come to the proof of inequality (2.14). The covariance function of the Ornstein–Uhlenbeck process starting from 0 writes

$$\Gamma^{\text{OU}}(s, t) = \frac{\sigma^2}{2\theta} e^{-\theta(s+t)} (e^{2\theta \min(s,t)} - 1).$$

For $t \in [0, T]$, we have $\int_0^T (\Gamma^{\text{OU}}(s, t))^2 \, ds = \frac{\sigma^4}{8\theta^3} (2 - 4e^{-2\theta t} \theta t - e^{-2\theta(T-t)} - 2e^{-2\theta t} + 2e^{-2\theta T} - e^{-2\theta(T+t)})$, and thus

$$\int_0^T \int_0^T (\Gamma^{\text{OU}}(s, t))^2 \, ds \, dt = \frac{\sigma^2}{16\theta^4} (-5 + 4\theta T + 8\theta T e^{-2\theta T} + 4e^{-2\theta T} + e^{-4\theta T}).$$

Consequently, the function ϕ defined by $\phi(\theta) := \int_0^T \int_0^T (\Gamma^{\text{OU}}(s, t))^2 \, ds \, dt - \frac{\sigma^4}{\theta^4}$ writes

$$\phi(\theta) = \frac{1}{16} \frac{\sigma^4}{\theta^4} (-21 + 4\theta T + 8\theta e^{-2\theta T} T + 4e^{-2\theta T} + e^{-4\theta T}).$$

Thus, $\phi(\theta) < -16 + 12\theta T$ which leads to inequality (2.14) thanks to the fact that $\theta T < \frac{4}{3}$.

We now come back to the general case where we might have $\theta T \geq \frac{4}{3}$. If this is the case, let us consider $\tilde{T} := T - \frac{1}{\theta}$, so that $\theta(T - \tilde{T}) < \frac{4}{3}$. For $t \in [\tilde{T}, T]$, we have

$$G_t = \underbrace{G_{\tilde{T}} e^{-\theta(t-\tilde{T})}}_{\text{independent of } (W_s)_{s \in [\tilde{T}, T]}} \stackrel{\perp}{+} \underbrace{\int_{\tilde{T}}^t \sigma e^{\theta(s-t)} \, dW_s}_{=: \tilde{G}_t^0}.$$

The so-defined process $(\tilde{G}_t^0)_{t \in [\tilde{T}, T]}$ is a centered Ornstein–Uhlenbeck process starting from 0 and satisfying the same SDE as G . Hence, by independence, if $\int_0^T g(s) \, dG_s = 0$, then $\int_{\tilde{T}}^T g(s) \, d\tilde{G}_s^0 = 0$.

As $\theta(T - \tilde{T}) < \frac{4}{3}$, we can apply the result to $(\tilde{G}_t^0)_{t \in [\tilde{T}, T]}$ so that $g|_{[\tilde{T}, T]} \stackrel{L^2([\tilde{T}, T])}{=} 0$. If $\tilde{T}\theta < \frac{4}{3}$, we then have $g \stackrel{L^2([0, T])}{=} 0$. If it is not the case, we use the same method by using the decomposition of $[0, \tilde{T}]$ into $[0, \tilde{T} - \frac{1}{\theta}]$ and $[\tilde{T} - \frac{1}{\theta}, \tilde{T}]$ and so on. An easy induction finally shows that $g \stackrel{L^2([0, T])}{=} 0$.

The converse is obvious. □

The case of a more general Gaussian semimartingale. In the [Appendix](#), we investigate the problem for more general Gaussian semimartingales. As we have seen in the case of the Ornstein–Uhlenbeck process, if functions $(f_i)_{i \in I}$ are linearly independent in $L^2([s, T], d\langle X \rangle)$ for $s \in [0, T]$, the (\mathcal{H}) hypothesis comes to the injectivity of the Wiener integral with respect to X on $\text{span}(f_i)_{i \in I}$ (on interval $[s, T]$).

3. K–L generalized bridges and partial functional quantization

We keep the notation and assumptions of Section 2.2. As we have seen, equation (2.3) decomposes the process X as the sum of a linear combination of the Karhunen–Loève coordinates $Y := (Y_i)_{i \in I}$ and an independent remainder term. We now consider \widehat{Y}^Γ a stationary Voronoi N -quantization of Y . \widehat{Y}^Γ can be written as a nearest neighbor projection of Y on a finite codebook $\Gamma = (\gamma_1, \dots, \gamma_N)$.

$$\widehat{Y}^\Gamma = \text{Proj}_\Gamma(Y) \quad \text{where } \text{Proj}_\Gamma \text{ is a nearest neighbor projection on } \Gamma.$$

For example, \widehat{Y}^Γ can be a stationary product quantization or an optimal quadratic quantization of Y . We now define the stochastic process $\widetilde{X}^{I, \Gamma}$ by replacing Y by \widehat{Y}^Γ in the decomposition (2.3). We denote $\widetilde{X}^{I, \Gamma} = \text{Proj}_{I, \Gamma}(X)$.

$$\widetilde{X}^{I, \Gamma} = \sum_{i \in I} \widehat{Y}_i^\Gamma e_i^X + \sum_{i \in \mathbb{N}^* \setminus I} \sqrt{\lambda_i^X \xi_i} e_i^X.$$

The conditional distribution of $\widetilde{X}^{I, \Gamma}$ given that Y falls in the Voronoi cell of γ_k is the probability distribution of the K–L generalized bridge with end-point γ_k . In other words, we have quantized the Karhunen–Loève coordinates of X corresponding to $i \in I$, and not the other ones.

The so-defined process $\widetilde{X}^{I, \Gamma}$ is called a *partial functional quantization of X* .

3.1. Partial functional quantization of stochastic differential equations

Let X be a continuous centered Gaussian semimartingale on $[0, T]$ with $X_0 = 0$. We consider the SDE

$$dS_t = b(t, S_t) dt + \sigma(t, S_t) dX_t, \quad S_0 = x \in \mathbb{R} \text{ and } t \in [0, T], \tag{3.1}$$

where $b(t, x)$ and $\sigma(t, x)$ are Borel functions, Lipschitz continuous with respect to x uniformly in t , σ and $b(\cdot, 0)$ are bounded. This SDE admits a unique strong solution S .

The conditional distribution given that $Y_i = y_i$ for $i \in I$ of S is the strong solution of the stochastic differential equation $dS_t = b(t, S_t) dt + \sigma(t, S_t) dX_t^{I, \bar{y}}$, with $S_0 = x \in \mathbb{R}$, and for $t \in [0, T]$, where $X_t^{I, \bar{y}}$ is the corresponding K–L generalized bridge.

Under the (\mathcal{H}) hypothesis, this suggests to define the partial quantization of S from a partial quantization $\tilde{X}^{I,\Gamma}$ of X by replacing X by $\tilde{X}^{I,\Gamma}$ in the SDE (3.1). We define the *partial quantization* $\tilde{S}^{I,\Gamma}$ as the process whose conditional distribution given that Y falls in the Voronoi cell of γ_k is the strong solution of the same SDE where X is replaced by the K–L generalized bridge with end-point γ_k . We write

$$d\tilde{S}_t^{I,\Gamma} = b(t, \tilde{S}_t^{I,\Gamma}) dt + \sigma(t, \tilde{S}_t^{I,\Gamma}) d\tilde{X}_t^{I,\Gamma}. \tag{3.2}$$

Remark. The SDE is written in the Itô sense unlike in the previous works on full functional quantization (Pagès and Sellami [26], Pagès and Printems [24]) where the SDE was written in the Stratonovich sense.

Here, the set I of quantized Karhunen–Loève coordinates does not depend on the quantization level, while in the case of full functional quantization, optimality is reached by adapting the quantization dimension. The optimal quantization dimension (or critical dimension) has been thoroughly investigated in Luschgy and Pagès [18,19] and is shown to be asymptotically equivalent to the logarithm of the quantization level when it goes to infinity, in the cases of Brownian motion, Brownian bridge and Ornstein–Uhlenbeck processes.

3.2. Convergence of partially quantized SDEs

We start by stating some useful inequalities for the sequel. Then we recall the so-called Zador’s theorem which will be used in the proof of the *a.s.* convergence of partially quantized SDEs.

Lemma 3.1 (Gronwall inequality for locally finite measures). *Consider \mathcal{I} an interval of the form $[a, b)$ or $[a, b]$ with $a < b$ or $[a, \infty)$. Let μ be a locally finite measure on the Borel σ -algebra of \mathcal{I} . We consider u a measurable function defined on \mathcal{I} such that for all $t \in \mathcal{I}$, $\int_a^t |u(s)|\mu(ds) < +\infty$. We assume that there exists a Borel function ψ on \mathcal{I} such that*

$$u(t) \leq \psi(t) + \int_{[a,t)} u(s)\mu(ds) \quad \forall t \in \mathcal{I}.$$

If $\begin{cases} \text{either } \psi \text{ is non-negative,} \\ \text{or } t \mapsto \mu([a, t)) \text{ is continuous on } \mathcal{I} \text{ and for all } t \in \mathcal{I}, \int_a^t |\psi(s)|\mu(ds) < \infty, \end{cases}$

then u satisfies the Gronwall inequality

$$u(t) \leq \psi(t) + \int_{[a,t)} \psi(s) \exp(\mu([s, t)))\mu(ds).$$

A proof of this result is available in Ethier and Kurtz [9], Appendix 5.1.

Lemma 3.2 (A Gronwall-like inequality in the non-decreasing case). *Consider \mathcal{I} an interval of the form $[a, b)$ or $[a, b]$ with $a < b$ or $[a, \infty)$. Let μ be a locally finite measure on the Borel*

σ -algebra of \mathcal{I} . We consider u a measurable non-decreasing function defined on \mathcal{I} such that for all $t \in \mathcal{I}$, $\int_a^t |u(s)|\mu(ds) < +\infty$. We assume that there exists a Borel function ψ on \mathcal{I} , and two non-negative constants $(A, B) \in \mathbb{R}_+^2$ such that

$$u(t) \leq \psi(t) + A \int_{[a,t)} u(s)\mu(ds) + B \sqrt{\int_{[a,t)} u(s)^2 \mu(ds)} \quad \forall t \in \mathcal{I}. \tag{3.3}$$

If $\begin{cases} \text{either } \psi \text{ is non-negative,} \\ \text{or } t \mapsto \mu([a, t)) \text{ is continuous on } \mathcal{I} \text{ and for all } t \in \mathcal{I}, \int_a^t |\psi(s)|\mu(ds) < \infty, \end{cases}$

then u satisfies the following Gronwall inequality

$$u(t) \leq 2\psi(t) + 2(2A + B^2) \int_{[a,t)} \psi(s) \exp((2A + B^2)\mu([s, t)))\mu(ds).$$

Proof. Using that for $(x, y) \in \mathbb{R}_+^2$, $\sqrt{xy} \leq \frac{1}{2}(\frac{x}{B} + By)$, we have

$$\left(\int_{[a,t)} u(s)^2 \mu(ds) \right)^{1/2} \leq \left(u(t) \int_{[a,t)} u(s)\mu(ds) \right)^{1/2} \leq \frac{u(t)}{2B} + \frac{B}{2} \int_{[a,t)} u(s)\mu(ds).$$

Plugging this in inequality (3.3) yields

$$u(t) \leq 2\psi(t) + (2A + B^2) \int_{[a,t)} u(s)\mu(ds).$$

Applying the regular Gronwall’s inequality (Lemma 3.1) yields the announced result. □

Theorem 3.1 (Zador, Bucklew, Wise, Graf, Luschgy, Pagès). Consider $r > 0$ and X be a \mathbb{R}^d -valued random variable such that $X \in L^{r+\eta}$ for some $\eta > 0$. We denote by $\mathcal{E}_{N,r}(X)$ the L^r optimal quantization error of level N for X , $\mathcal{E}_{N,r}(X) := \min\{\|X - Y\|_r, |Y(\Omega)| \leq N\}$.

1. (Sharp rate). Let $\mathbb{P}_X(d\xi) = \phi(\xi) d\xi + \nu(d\xi)$ be the Radon–Nikodym decomposition of the probability distribution of X . (ν and the Lebesgue’s measure are singular). Then if $\phi \not\equiv 0$,

$$\mathcal{E}_{N,r}(X) \underset{N \rightarrow \infty}{\sim} \tilde{J}_{r,d} \times \left(\int_{\mathbb{R}^d} \phi^{d/(d+r)}(u) du \right)^{1/d+1/r} \times N^{-1/d},$$

where $\tilde{J}_{r,d} \in (0, \infty)$.

2. (Non-asymptotic upper bound). There exists $C_{d,r,\eta} \in (0, \infty)$ such that, for every \mathbb{R}^d -valued random vector X ,

$$\forall N \geq 1, \quad \mathcal{E}_{N,r}(X) \leq C_{d,r,\eta} \|X\|_{r+\eta} N^{-1/d}.$$

The first statement of the theorem was first established for probability distributions with compact support by Zador [30], and extended by Bucklew and Wise to general probability distributions on \mathbb{R}^d (Bucklew and Wise [4]). The first mathematically rigorous proof can be found in Graf and Luschgy [11]. The proof of the second statement is available in Luschgy and Pagès [20].

The real constant $J_{r,d}$ corresponds to the case of the uniform probability distribution over the unit hypercube $[0, 1]^d$. We have $\tilde{J}_{r,1} = \frac{1}{2}(r + 1)^{-1/r}$ and $\tilde{J}_{2,2} = \sqrt{\frac{5}{18\sqrt{2}}}$ (see Graf and Luschgy [11]).

3.2.1. L^p convergence of partially quantized SDEs

Lemma 3.3 (Generalized Minkowski inequality for locally finite measures). *Consider \mathcal{I} an interval of the form $[a, b)$ or $[a, b]$ with $a < b$ or $[a, \infty)$. Let μ be a locally finite measure on the Borel σ -algebra of \mathcal{I} . Then for any non-negative bi-measurable process $X = (X_t)_{t \in \mathcal{I}}$ and every $p \in [1, \infty)$,*

$$\left\| \int_{\mathcal{I}} X_t \mu(dt) \right\|_p \leq \int_{\mathcal{I}} \|X_t\|_p \mu(dt).$$

Proposition 3.1 (Burkholder–Davis–Gundy inequality). *For every $p \in (0, \infty)$, there exist two positive real constants c_p^{BDG} and C_p^{BDG} such that for every continuous local martingale $(X_t)_{t \in [0, T]}$ null at 0,*

$$c_p^{\text{BDG}} \|\sqrt{\langle X \rangle_T}\|_p \leq \left\| \sup_{s \in [0, T]} |X_s| \right\|_p \leq C_p^{\text{BDG}} \|\sqrt{\langle X \rangle_T}\|_p.$$

We refer to Revuz and Yor [27] for a detailed proof.

Proposition 3.2 (L^p inequality). *Let G be a standard Gaussian random variable valued in \mathbb{R} . There exists a constant $C_p > 0$ such that for every $M > 1$*

$$\begin{aligned} & \sqrt{\frac{2}{\pi}} M^{p-1} \exp\left(-\frac{M^2}{2}\right) \\ & \leq \mathbb{E}[|G|^p \mathbf{1}_{|G|>M}] \leq C_p M^{p-1} \exp\left(-\frac{M^2}{2}\right). \end{aligned}$$

Consequently

$$\begin{aligned} & \left(\sqrt{\frac{2}{\pi}}\right)^{1/p} M^{1/q} \exp\left(-\frac{M^2}{2p}\right) \\ & \leq \|G \mathbf{1}_{|G|>M}\|_p \leq (C_p)^{1/p} M^{1/q} \exp\left(-\frac{M^2}{2p}\right), \end{aligned}$$

where q is the conjugate exponent of p .

Proposition 3.3 (The non-standard case and L^p reverse inequality). *If $H := \sigma G$ has a variance of σ^2 , we obtain*

$$\begin{aligned} \|H\mathbf{1}_{|H|>M}\|_p &\leq \sigma \|G\mathbf{1}_{|G|>M/\sigma}\|_p = \sigma (C_p)^{1/p} \left(\frac{M}{\sigma}\right)^{1/q} \exp\left(-\frac{M^2}{2p\sigma^2}\right) \\ &= \underbrace{\sigma^{1/p} (C_p)^{1/p} M^{1/q}}_{=: \eta_M} \exp\left(-\frac{M^2}{2p\sigma^2}\right). \end{aligned} \tag{3.4}$$

Conversely, for some fixed $\eta > 0$, and if $M > 1$, we have

$$M \geq \underbrace{\sqrt{-\sigma^2(p-1)\mathcal{W}_{-1}\left(-\frac{q\eta^{2q}}{p\sigma^2(C_p^{2q/p}\sigma^{2q/p})}\right)}}_{=: M_\eta} \Rightarrow \eta_M \leq \eta, \tag{3.5}$$

where \mathcal{W}_{-1} is the secondary branch of the Lambert \mathcal{W} function. For more details on the Lambert \mathcal{W} function, we refer to Corless et al. [7].

Theorem 3.2 (L^p quantization of partially quantized SDEs). *Let X be a continuous centered Gaussian martingale on $[0, T]$ with $X_0 = 0$. Let S be the strong solution of the SDE*

$$dS_t = b(t, S_t) dt + \sigma(t, S_t) dX_t, \quad S_0 = x,$$

where $b(t, x)$ and $\sigma(t, x)$ are Borel functions, Lipschitz continuous with respect to x uniformly in t , σ and $b(\cdot, 0)$ are bounded.

We consider $\tilde{X}^{I,\Gamma}$ a stationary partial functional quantization of X and $\tilde{S}^{I,\Gamma}$ the corresponding partial functional quantization of S , that is, the strong solutions of

$$d\tilde{S}_t^{I,\Gamma} = b(t, \tilde{S}_t^{I,\Gamma}) dt + \sigma(t, \tilde{S}_t^{I,\Gamma}) d\tilde{X}_t^{I,\Gamma}, \quad \tilde{S}_0^{I,\Gamma} = x.$$

Then, for every $p \in (0, \infty)$, $\varepsilon > 0$ and $t \in [0, T]$, there exist three positive constants $C_{X,\varepsilon,I}$, $A_{X,\varepsilon,I}$ and $B_{X,\varepsilon,I}$ such that

$$\left\| \sup_{v \in [0,t]} |S_v - \tilde{S}_v^{I,\Gamma}| \right\|_p \leq C_{X,\varepsilon,I} \exp\left(A_{X,\varepsilon,I} \sqrt{-\mathcal{W}_{-1}\left(-\frac{\|Y - \hat{Y}^\Gamma\|_{p+\varepsilon}^{2q}}{B_{X,\varepsilon,I}}\right)} \right) \|Y - \hat{Y}^\Gamma\|_{p+\varepsilon}, \tag{3.6}$$

where q is the conjugate exponent of p , where Y is defined from X by equation (2.3) and \hat{Y}^Γ is the nearest neighbor projection on Γ .

Remark. Using that $\mathcal{W}_{-1}(-x) \underset{x \rightarrow 0_+}{\sim} \ln(x)$, we can see that the right-hand term in equation (3.6) goes to 0 as the quantization error $\|Y - \hat{Y}^\Gamma\|_{p+\varepsilon}$ goes to 0.

Proof. We decompose the process X into $X_t = \sum_{i \in I} Y_i e_i^X(t) + X_t^{I, \bar{0}}$ and $\tilde{X}^{I, \Gamma}$ into $\tilde{X}_t^{I, \Gamma} = \sum_{i \in I} \hat{Y}_i^\Gamma e_i^X(t) + X_t^{I, \bar{0}}$, where \hat{Y}^Γ is the nearest neighbor projection of Y on Γ . For some $k \in \{1, \dots, N\}$, conditionally to $\hat{Y}^\Gamma = \gamma_k$, we have

$$\begin{aligned} S_t - \tilde{S}_t^{I, \Gamma} &= \int_0^t (b(u, S_u) - b(u, \tilde{S}_u^{I, \Gamma})) du + \sum_{i \in I} \int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) \hat{Y}_i^\Gamma de_i^X(u) \\ &\quad + \sum_{i \in I} \int_0^t (Y_i - \hat{Y}_i^\Gamma) \sigma(u, S_u) de_i^X(u) + \int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) G_u d\langle X \rangle_u \\ &\quad + \int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) d\tilde{M}_u. \end{aligned}$$

This gives (conditionally to $\hat{Y}^\Gamma = \gamma_k$)

$$\begin{aligned} &|S_t - \tilde{S}_t^{I, \Gamma}| \\ &\leq [b]_{\text{Lip}} \int_0^t |S_u - \tilde{S}_u^{I, \Gamma}| du + [\sigma]_{\text{Lip}} |I| \max_{u \in [0, T]} |(e_i^X)'(u)| \left(\max_{i \in I} |\hat{Y}_i^\Gamma| \right) \int_0^t |S_u - \tilde{S}_u^{I, \Gamma}| du \\ &\quad + [\sigma]_{\text{max}} |I| \max_{u \in [0, T]} |(e_i^X)'(u)| T \sum_{i \in I} |Y_i - \hat{Y}_i^\Gamma| + \left| \int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) G_u d\langle X \rangle_u \right| \\ &\quad + \left| \int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) d\tilde{M}_u \right|. \end{aligned}$$

As a consequence, conditionally to $\hat{Y}^\Gamma = \gamma_k$,

$$\begin{aligned} \max_{v \in [0, t]} |S_v - \tilde{S}_v^{I, \Gamma}| &\leq [b]_{\text{Lip}} \int_0^t \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| du \\ &\quad + [\sigma]_{\text{Lip}} |I| \max_{u \in [0, T]} |(e_i^X)'(u)| \left(\max_{i \in I} |\hat{Y}_i^\Gamma| \right) \int_0^t \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| du \\ &\quad + [\sigma]_{\text{max}} |I| \max_{u \in [0, T]} |(e_i^X)'(u)| T \sum_{i \in I} |Y_i - \hat{Y}_i^\Gamma| \\ &\quad + \max_{v \in [0, t]} \left| \int_0^v (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) G_u d\langle X \rangle_u \right| \\ &\quad + \max_{v \in [0, t]} \left| \int_0^v (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) d\tilde{M}_u \right|. \end{aligned}$$

To shorten the notation, we denote, for a random variable V and a non-negligible event A , $\|V\|_{p, A} := \mathbb{E}[V^p | A]^{1/p}$. Hence, using the Minkowski inequality and the generalized Minkowski

inequality for locally finite measures (Lemma 3.3), we get

$$\begin{aligned}
 & \left\| \max_{v \in [0, t]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \leq [b]_{\text{Lip}} \int_0^t \left\| \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} du \\
 & \quad + [\sigma]_{\text{Lip}} |I| \max_{\substack{i \in I \\ u \in [0, T]}} |(e_i^X)'(u)| \left(\max_{i \in I} |\widehat{Y}_i^\Gamma| \right) \int_0^t \left\| \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} du \\
 & \quad + [\sigma]_{\text{Lip}} |I| \max_{\substack{i \in I \\ u \in [0, T]}} |(e_i^X)'(u)| T \left\| \sum_{i \in I} |Y_i - \widehat{Y}_i^\Gamma| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \quad + \left\| \max_{v \in [0, t]} \left| \int_0^v (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) G_u d\langle X \rangle_u \right| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \quad + \left\| \max_{v \in [0, t]} \left| \int_0^v (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})) d\tilde{M}_u \right| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}}.
 \end{aligned}$$

Now, from the Burkholder–Davis–Gundy inequality,

$$\begin{aligned}
 & \left\| \max_{v \in [0, t]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \leq [b]_{\text{Lip}} \int_0^t \left\| \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} du \\
 & \quad + [\sigma]_{\text{Lip}} |I| \max_{\substack{i \in I \\ u \in [0, T]}} |(e_i^X)'(u)| \left(\max_{i \in I} |\widehat{Y}_i^\Gamma| \right) \int_0^t \left\| \max_{v \in [0, u]} |S_v - \tilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} du \\
 & \quad + [\sigma]_{\text{Lip}} |I| \max_{\substack{i \in I \\ u \in [0, T]}} |(e_i^X)'(u)| T \left\| \sum_{i \in I} |Y_i - \widehat{Y}_i^\Gamma| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \quad + \left\| \int_0^t |\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma})| |G_u| d\langle X \rangle_u \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \quad + C_p^{\text{BDG}} \left\| \sqrt{\int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I, \Gamma}))^2 d\langle X \rangle_u} \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}}.
 \end{aligned} \tag{3.7}$$

Now, from Schwarz’s inequality

$$\left\| \sum_{i \in I} |Y_i - \widehat{Y}_i^\Gamma| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \leq \left\| \sqrt{|I|} \sqrt{\sum_{i \in I} |Y_i - \widehat{Y}_i^\Gamma|^2} \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} = \sqrt{|I|} \|Y - \widehat{Y}^\Gamma\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}}.$$

From the generalized Minkowski inequality

$$\begin{aligned}
 & \left\| \int_0^t |\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})| |G_u| d\langle X \rangle_u \right\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} \\
 & \leq \int_0^t \|(\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})) G_u\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u \\
 & = \int_0^t \|(\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})) G_u \mathbf{1}_{|G_u| \geq M} \\
 & \quad + (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})) G_u \mathbf{1}_{|G_u| \leq M}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u \\
 & \leq \int_0^t \|(\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})) G_u \mathbf{1}_{|G_u| \geq M}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u \\
 & \quad + \int_0^t \|(\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})) G_u \mathbf{1}_{|G_u| \leq M}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u \\
 & \leq 2[\sigma]_{\max} \int_0^t \|G_u \mathbf{1}_{|G_u| \geq M}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u + M[\sigma]_{\text{Lip}} \int_0^t \|S_u - \tilde{S}_u^{I,\Gamma}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u.
 \end{aligned}$$

We obtain, thanks to Proposition 3.3

$$\begin{aligned}
 & \left\| \int_0^t |\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma})| |G_u| d\langle X \rangle_u \right\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} \\
 & \leq \underbrace{2[\sigma]_{\max} \langle X \rangle_t (C_p)^{1/p} v_t^{1/p} M^{1/q} \exp\left(-\frac{M^2}{2pv_t^2}\right)}_{=: \eta_M} + M[\sigma]_{\text{Lip}} \int_0^t \|S_u - \tilde{S}_u^{I,\Gamma}\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u,
 \end{aligned}$$

where $v_t^2 = \max_{u \in [0,t]} (\text{Var}(G_u))$. Moreover,

$$\left\| \sqrt{\int_0^t (\sigma(u, S_u) - \sigma(u, \tilde{S}_u^{I,\Gamma}))^2 d\langle X \rangle_u} \right\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} \leq \sqrt{\int_0^t \max_{\substack{i \in I \\ v \in [0,u]}} |S_v - \tilde{S}_v^{I,\Gamma}|^2 d\langle X \rangle_u}.$$

Hence, equation (3.7) becomes

$$\begin{aligned}
 & \left\| \max_{v \in [0,t]} |S_v - \tilde{S}_v^{I,\Gamma}| \right\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} \\
 & \leq \underbrace{[\sigma]_{\text{Lip}} |I| \max_{\substack{i \in I \\ u \in [0,T]}} |(e_i^X)'(u)| \sqrt{|I|}}_{=: A_i^X} \|Y - \hat{Y}^\Gamma\|_{p, \{\hat{Y}^\Gamma = \gamma_k\}} + \eta_M \\
 & \quad + [b]_{\text{Lip}} \int_0^t \max_{v \in [0,u]} |S_v - \tilde{S}_v^{I,\Gamma}| d\langle X \rangle_u
 \end{aligned}$$

$$\begin{aligned}
 & + [\sigma]_{\text{Lip}} |I| \max_{i \in I} |(e_i^X)'(u)| \left(\max_{i \in I} |\widehat{Y}_i^\Gamma| \right) \int_0^t \left\| \max_{v \in [0, u]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} du \\
 & + C_p^{\text{BDG}} \left(\int_0^t 2 \left\| \max_{i \in I} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}}^2 d\langle X \rangle_u \right)^{1/2} \\
 & + \underbrace{M[\sigma]_{\text{Lip}}}_{=: C^{X, M}} \int_0^t \left\| \max_{v \in [0, u]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} d\langle X \rangle_u.
 \end{aligned}$$

We can then apply the ‘‘Gronwall-like’’ Lemma 3.2 for locally finite measures to the non-decreasing function

$$\left\| \sup_{v \in [0, t]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} = \mathbb{E} \left[\sup_{v \in [0, t]} |S_v - \widetilde{S}_v^{I, \Gamma}|^p | \widehat{Y}^\Gamma = \gamma_k \right]^{1/p}$$

and with the locally finite measure μ defined by $\mu(du) = du + d\langle X \rangle_u$, and we obtain

$$\begin{aligned}
 & \left\| \sup_{v \in [0, t]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_{p, \{\widehat{Y}^\Gamma = \gamma_k\}} \\
 & \leq (A_I^X \mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma = \gamma_k]^{1/p} + \eta_M) \exp((E_I^{X, \gamma_k} + C^{X, M})\mu([0, t])) \\
 & \leq (A_I^X \mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma = \gamma_k]^{1/p} + \eta_M) \underbrace{\exp(E_I^{X, \gamma_k} \mu([0, t]))}_{=: \phi(\gamma_k)} \exp(C^{X, M} \mu([0, t])),
 \end{aligned}$$

where E_I^{X, γ_k} is an affine function of $\max_{i \in I} |(\gamma_k)_i|$. This yields

$$\begin{aligned}
 & \left\| \sup_{v \in [0, t]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_p \\
 & \leq (A_I^X \|\mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma]^{1/p} \phi(\widehat{Y}^\Gamma)\|_p + \eta_M \|\phi(\widehat{Y}^\Gamma)\|_p) \exp(C^{X, M} \mu([0, t])).
 \end{aligned}$$

Now, for $\varepsilon > 0$ and $\tilde{p} = 1 + \frac{\varepsilon}{p}$ and $\tilde{q} = \frac{\tilde{p}}{\tilde{p}-1} = 1 + \frac{p}{\varepsilon}$ the conjugate exponent of \tilde{p} , we have, thanks to Hölder’s inequality

$$\begin{aligned}
 \mathbb{E}[\phi(\widehat{Y}^\Gamma)^p \mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma]] & \leq \|\phi(\widehat{Y}^\Gamma)^p\|_{\tilde{q}} \|\mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma]\|_{\tilde{p}} \\
 & \leq \|\phi(\widehat{Y}^\Gamma)^p\|_{\tilde{q}} \mathbb{E}[|Y - \widehat{Y}^\Gamma|^{p+\varepsilon}]^{p/(p+\varepsilon)}.
 \end{aligned}$$

Hence,

$$\|\mathbb{E}[|Y - \widehat{Y}^\Gamma|^p | \widehat{Y}^\Gamma]^{1/p} \phi(\widehat{Y}^\Gamma)\|_p \leq \|\phi(\widehat{Y}^\Gamma)^p\|_{\tilde{q}}^{1/p} \mathbb{E}[|Y - \widehat{Y}^\Gamma|^{p+\varepsilon}]^{1/(p+\varepsilon)}.$$

Now, as the so-defined function ϕ is convex and as \widehat{Y}^Γ is a stationary quantizer of Y , we have thanks to equation (1.4), $\|\phi(\widehat{Y}^\Gamma)^p\|_{\tilde{q}} \leq \|\phi(Y)^p\|_{\tilde{q}}$ and $\|\phi(\widehat{Y}^\Gamma)\|_p \leq \|\phi(Y)\|_p$.

If one sets

$$M = \sqrt{-v_t(p-1)\mathcal{W}_{-1}\left(-\frac{q\|Y - \widehat{Y}^\Gamma\|_{p+\varepsilon}^{2q}}{pv_t^2 C_p^{2q/p} v_t^{2q/p}}\right)},$$

where q is the conjugate exponent of p and \mathcal{W}_{-1} is the secondary branch of the Lambert \mathcal{W} function, Proposition 3.3 ensures that $\eta_M \leq \eta := \|Y - \widehat{Y}^\Gamma\|_{p+\varepsilon}$. We finally have the following error bound

$$\begin{aligned} & \left\| \sup_{v \in [0, t]} |S_v - \widetilde{S}_v^{I, \Gamma}| \right\|_p \\ & \leq C_{X, \varepsilon, I} \exp\left([\sigma]_{\text{Lip}} \sqrt{-v_t(p-1)\mathcal{W}_{-1}\left(-\frac{q\|Y - \widehat{Y}^\Gamma\|_{p+\varepsilon}^{2q}}{pv_t^2 C_p^{2q/p} v_t^{2q/p}}\right)}\right) \|Y - \widehat{Y}^\Gamma\|_{p+\varepsilon}, \end{aligned}$$

which is the desired inequality. □

Remark (Without the stationarity property). The last step of the proof of Theorem 3.2 (the use of Jensen’s inequality) relies on the stationarity of the quantizer \widehat{Y} . Now, without this stationarity hypothesis and under the additional assumption that

$$\Gamma \cap B(0, 1) \neq \emptyset, \tag{A}$$

we have for every $i \in I$

$$|\widehat{Y}_i| \leq |Y_i - \widehat{Y}_i| + |Y_i| \leq |Y_i| + |Y_i - \gamma_i^{k_0}| \leq 2|Y_i| + |\gamma_i^{k_0}| \leq 2|Y_i| + 1 \quad \text{where } \gamma^{k_0} \in \Gamma \cap B(0, 1).$$

Hence,

$$\max_{i \in I} |\widehat{Y}_i| \leq 2 \max_{i \in I} |Y_i| + 1.$$

We notice that the function $\phi(x)$ defined in the demonstration of Theorem 3.2 writes $\phi(x) = \psi(\max_{i \in I} x_i)$ for some non-decreasing function ψ . This implies

$$\phi(\widehat{Y}) = \psi\left(\max_{i \in I} \widehat{Y}_i\right) \leq \psi\left(\max_{i \in I} (2|Y_i| + 1)\right) = \phi(2|Y| + 1).$$

Hence, we can obtain the same conclusion as in Theorem 3.2.

Corollary 3.1 (L^p convergence). *With the same notation and hypothesis as in Theorem 3.2, consider $(\widetilde{X}^{I, \Gamma_n})_{n \in \mathbb{N}}$ a sequence of partial functional quantizers of X and $(\widetilde{S}^{I, \Gamma_n})_{n \in \mathbb{N}}$ the corresponding sequence of partial quantizers of S . (For $n \in \mathbb{N}$, Γ_n is assumed to have cardinal n .)*

If we make the additional assumption that the associated sequence of quantizers $(\widehat{Y}^{\Gamma_n})_{n \in \mathbb{N}}$ is rate-optimal for the $L^{p+\varepsilon}$ convergence for some $\varepsilon > 0$, then for every $t \in [0, T)$ we have

$$\mathbb{E}\left[\sup_{u \in [0, t]} |S_u - \widetilde{S}_u^{I, \Gamma_n}|^p\right]^{1/p} = \mathcal{O}(n^{-1/|I|}).$$

Proof. As $\|Y - \widehat{Y}^{\Gamma_n}\|_{p+\varepsilon} \rightarrow_{n \rightarrow \infty} 0$, we have a.s. $d(\widehat{Y}^{\Gamma_n}, Y) \rightarrow_{n \rightarrow \infty} 0$. Hence, there exists $N_0 \in \mathbb{N}$ such that for every $n \geq N_0$, Γ_n satisfies Hypothesis (A). From this observation, the result is straightforward consequence of the previous remark and Zador’s theorem 3.1, which defines the optimal convergence rate of a sequence of quantizers. \square

3.2.2. *The a.s. convergence of partially quantized SDEs*

Theorem 3.3 (Almost sure convergence of partially quantized SDEs). *Let X be a continuous centered Gaussian martingale on $[0, T]$ with $X_0 = 0$. Let S be the strong solution of the SDE*

$$dS_t = b(t, S_t) dt + \sigma(t, S_t) dX_t, \quad S_0 = x,$$

where $b(t, x)$ and $\sigma(t, x)$ are Borel functions, Lipschitz continuous with respect to x uniformly in t , σ and $b(\cdot, 0)$ are bounded.

We consider $(\widetilde{X}^{I, \Gamma_k})_{k \in \mathbb{N}}$ a sequence of partial functional quantizers of X and $\widetilde{S}^{I, \Gamma_n}$ the corresponding partial functional quantization of S , that is, the strong solutions of

$$d\widetilde{S}_t^{I, \Gamma_n} = b(t, \widetilde{S}_t^{I, \Gamma_n}) dt + \sigma(t, \widetilde{S}_t^{I, \Gamma_n}) d\widetilde{X}_t^{I, \Gamma_n}, \quad \widetilde{S}_0^{I, \Gamma_n} = x.$$

(For $n \in \mathbb{N}$, Γ_n is assumed to have cardinal n .) We also assume that the sequence of partial quantizers of X is rate-optimal for some $p > |I|$, that is, that there exists a constant C such that

$$\mathbb{E}[|Y - \widehat{Y}^{\Gamma_n}|^p]^{1/p} \leq Cn^{-1/|I|}$$

for every $n \in \mathbb{N}^*$, where Y is defined from X by equation (2.3) and \widehat{Y}^Γ is the nearest neighbor projection on Γ . Then for every $t \in [0, T]$, $\widetilde{S}_t^{I, \Gamma_n}$ converges almost surely to S_t .

Proof. From Corollary 3.1, if $t \in [0, T]$, there exist $r \in (|I|, p)$ and $N_0 \in \mathbb{N}$ such that for $n \geq N_0$,

$$\mathbb{E}\left[\sup_{u \in [0, t]} |S_u - \widetilde{S}_u^{I, \Gamma_n}|^r\right]^{1/r} = O(n^{-1/|I|}).$$

Hence, as $\frac{r}{|I|} > 1$, Beppo–Levi’s theorem for series with non-negative terms implies

$$\mathbb{E}\left[\sum_{n \geq 1} \sup_{u \in [0, t]} |S_u - \widetilde{S}_u^{I, \Gamma_n}|^r\right] < +\infty.$$

Thus $\sum_{n \geq 1} \sup_{u \in [0, t]} |S_u - \widetilde{S}_u^{I, \Gamma_n}|^r < +\infty$ \mathbb{P} -a.s. so that $\sup_{u \in [0, t]} |S_u - \widetilde{S}_u^{I, \Gamma_n}| \rightarrow_{n \rightarrow \infty} 0$ \mathbb{P} -a.s. \square

Remark (Extension to semimartingales). In Theorems 3.2 and 3.3, we limited ourselves to the case where X is a local martingale. The proofs are easily extended to the case of a semimartingale X as soon as there exists a locally finite measure ν on $[0, T]$ such that for every $\omega \in \Omega$ the finite-variation part $dV(\omega)$ in the canonical decomposition of X is absolutely continuous with respect to ν . In particular, this is the case for the Brownian bridge and Ornstein–Uhlenbeck processes whose finite-variation parts are absolutely continuous with respect to the Lebesgue measure on $[0, T]$.

Appendix: Injectivity properties of the Wiener integral

In this Appendix, we recall some results on the definition of the Wiener integral with respect to a Gaussian process. We focus on the injectivity properties. Here, we pay special attention to the special case of the Ornstein–Uhlenbeck processes.

The covariance operator and the Cameron–Martin space

Consider X a bi-measurable centered Gaussian process on $[0, T]$ such that $\int_0^T \mathbb{E}[X_t^2] dt < \infty$ and with a continuous covariance function Γ^X on $[0, T] \times [0, T]$. We denote by $H := \overline{\text{span}\{X_t, t \in [0, T]\}}^{L^2(\mathbb{P})}$ the Gaussian Hilbert space spanned by $(X_t)_{t \in [0, T]}$. The covariance operator C_X of X is defined by

$$C_X : L^2([0, T]) \rightarrow L^2([0, T]),$$

$$y \mapsto C_X y = \mathbb{E}[(y, X)X].$$

We have $C_X y(t) = \mathbb{E}[(y, X)X](t) = \mathbb{E}[\int_0^T X_s y(s) ds X_t] = \int_0^T \Gamma^X(t, s) y(s) ds$ where $\Gamma^X(t, s) = \mathbb{E}[X_t X_s]$ is the covariance function of X .

The Cameron–Martin space of X , (or reproducing Hilbert space of C_X), which we denote by K_X , is the subspace of $L^2([0, T])$ defined by $K_X := \{t \mapsto \mathbb{E}[Z X_t], Z \in H\}$. K_X is equipped with the scalar product defined by

$$\langle k_1, k_2 \rangle_X = \mathbb{E}[Z_1 Z_2] \quad \text{if } k_i = \mathbb{E}[Z_i X \cdot], i = 1, 2,$$

so that $(K_X, \langle \cdot, \cdot \rangle_X)$ is a Hilbert space, isometric with the Hilbert space $\overline{\{(y, X) : y \in L^2([0, T])\}}^H$. K_X is spanned as a Hilbert space by $\{C_X(y) : y \in L^2([0, T])\}$.

The Wiener integral

Here, we follow the same steps as in Lebovits and Lévy–Véhel [17] and in Jost [16] for the definition of a general Wiener integral. The difference here is that we use the quotient topology in order to define the Wiener integral in a more general setting.

We define the map $U : H \rightarrow K_X$ defined by $U(Z)(t) = \mathbb{E}[Z X_t]$. By definition of H and K_X , U is a bijection and for any $s \in [0, T]$, we have $U(X_s) = \Gamma^X(s, \cdot)$. Consequently, K_X is spanned by $(\Gamma^X(s, \cdot))_{s \in [0, T]}$ as a Hilbert space. Now, we linearly map the set of the piecewise constant functions $\mathcal{E}([0, T])$ to the Cameron–Martin space K_X by

$$J : \mathcal{E}([0, T]) \rightarrow K_X,$$

$$\mathbf{1}_{|s, t|} \mapsto \Gamma^X(t, \cdot) - \Gamma^X(s, \cdot),$$

where $|a, b|$ stands either for the interval $[a, b]$, (a, b) , $(a, b]$ or $[a, b)$. We equip $\mathcal{E}([0, T])$ with the bilinear form $\langle \cdot, \cdot \rangle_J$ which is defined by

$$\langle f, g \rangle_J := \langle Jf, Jg \rangle_X.$$

It is a bilinear symmetric positive-semidefinite form.

Remark. The so-called reproducing property shows that $\langle \mathbf{1}_{[0,t]}, \mathbf{1}_{[0,s]} \rangle_J = \Gamma^X(t, s) + \Gamma^X(0, 0) - \Gamma^X(0, s) - \Gamma^X(0, t)$. When $X_0 = 0$ a.s., this gives $\langle \mathbf{1}_{[0,t]}, \mathbf{1}_{[0,s]} \rangle_J = \Gamma^X(s, t)$.

Now, we define the equivalence relation \sim_J on $\mathcal{E}([0, T])$ by $x \sim_J y$ if $\langle x - y, x - y \rangle_J = 0$. On the quotient space $E([0, T]) := \mathcal{E}([0, T]) / \sim_J$, the bilinear form $\langle \cdot, \cdot \rangle_J$ is positive-definite and thus it is a scalar product on $E([0, T])$. In this context, J defines an (isometric) linear map from $E([0, T])$ to K_X . Then, considering the completion F of $E([0, T])$ associated with this scalar product, J is extended to F and $U^{-1} \circ J : F \rightarrow H$ is an (isometric) injective map that we call Wiener integral associated to X .

$$\int_0^T f(t) dX_t := U^{-1} \circ J(f).$$

Injectivity properties of the Wiener integral

As we have just seen, the Wiener integral is an (isometric) injective map from F to H . Still, for example, when dealing with a standard Brownian bridge on $[0, T]$, $\|\mathbf{1}_{[0,T]}\|_J = 0$, so that there are functions of $\mathcal{E}([0, T])$ which have a non-zero L^2 norm and a zero $\|\cdot\|_J$ norm. Injectivity only holds in the quotient space $E([0, T]) = \mathcal{E}([0, T]) / \sim_J$ and its completion F .

It is classical background that in the special case of a standard Brownian motion, $\|\cdot\|_J$ exactly coincides with the canonical L^2 norm so that $F = L^2([0, T])$.

Study of the case of Ornstein–Uhlenbeck processes

From now, we will assume that X is a centered Ornstein–Uhlenbeck process defined on $[0, T]$ by the SDE

$$dX_t = -\theta X_t dt + \sigma dW_t \quad \text{with } \sigma > 0 \text{ and } \theta > 0,$$

where W is a standard Brownian motion and $X_0 \stackrel{\mathcal{L}}{\sim} \mathcal{N}(0, \sigma_0^2)$ is independent of W . We make the additional assumption that $\theta T \leq \frac{4}{3}$. The covariance function writes

$$\Gamma^X(s, t) = \frac{\sigma^2}{2\theta} e^{-\theta(s+t)} (e^{2\min(s,t)} - 1) + \sigma_0^2 e^{-\theta(s+t)}.$$

Proposition A.1 (Semi-norm equivalence on $\mathcal{E}([0, T])$). *There exist two positive constants c and C such that for every $f \in \mathcal{E}([0, T])$, $c\|f\|_2 \leq \|f\|_J \leq C\|f\|_2$.*

Proof. Let us consider $f \in \mathcal{E}([0, T])$. We have

$$\begin{aligned} \|f\|_J^2 &= \text{Var} \left(-\theta \int_0^T f(s) X_s ds + \sigma \int_0^T f(s) dW_s \right) \\ &\leq 2 \text{Var} \left(\theta \int_0^T f(s) X_s ds \right) + 2 \text{Var} \left(\sigma \int_0^T f(s) dW_s \right). \end{aligned}$$

The solution of the Ornstein–Uhlenbeck SDE is

$$X_t = X_0 e^{-\theta t} + \underbrace{\int_0^t \sigma e^{\theta(s-t)} dW_s}_{=: X_t^0}.$$

The so-defined process $(X_t^0)_{t \in [0, T]}$ is a centered Ornstein–Uhlenbeck process starting from 0. Hence, we have

$$\begin{aligned} \|f\|_J^2 &\leq 2 \operatorname{Var}\left(X_0 \theta \int_0^T f(s) e^{-\theta s} ds\right) + 2 \operatorname{Var}\left(\theta \int_0^T f(s) X_s^0 ds\right) + 2 \operatorname{Var}\left(\sigma \int_0^T f(s) dW_s\right) \\ &\leq 2\theta^2 T \operatorname{Var}(X_0) \int_0^T f(s)^2 ds + 2 \operatorname{Var}\left(\theta \int_0^T f(s) X_s^0 ds\right) + 2 \operatorname{Var}\left(\sigma \int_0^T f(s) dW_s\right). \end{aligned}$$

As in the proof of Proposition 2.6, using that $\theta T < 4/3$, we can show that $\operatorname{Var}(\theta \int_0^T f(s) X_s^0 ds) \leq \operatorname{Var}(\sigma \int_0^T f(s) dW_s)$. Hence,

$$\|f\|_J^2 \leq \underbrace{(2\theta^2 T \sigma_0^2 + 4\sigma^2)}_{=: C^2} \int_0^T f(s)^2 ds,$$

which is the desired inequality. Now we write

$$\int_0^t f(s) dX_s = \underbrace{-\theta \int_0^t f(s) X_0 e^{-\theta s} ds}_{=: G_0^f} + \underbrace{\left(-\theta \int_0^t f(s) X_s^0 ds\right)}_{=: G_1^f} + \underbrace{\sigma \int_0^t f(s) dW_s}_{=: G_2^f},$$

where (G_0^f, G_1^f, G_2^f) is Gaussian and G_0^f is independent of G_1^f and G_2^f . Hence,

$$\begin{aligned} \operatorname{Var}\left(\int_0^t f(s) dX_s\right) &\geq \operatorname{Var}(G_1^f + G_2^f) = \operatorname{Var}(G_1^f) + \operatorname{Var}(G_2^f) + 2 \operatorname{cov}(G_1^f, G_2^f) \\ &\geq \operatorname{Var}(G_1^f) + \operatorname{Var}(G_2^f) - 2\sqrt{\operatorname{Var}(G_1^f) \operatorname{Var}(G_2^f)} = \left(\sqrt{\operatorname{Var}(G_2^f)} - \sqrt{\operatorname{Var}(G_1^f)}\right)^2. \end{aligned} \tag{A.1}$$

It has been shown at the beginning of the proof of Proposition 2.6 that there exists a constant $K < 1$ independent of f such that $\operatorname{Var}(G_1^f) \leq K \operatorname{Var}(G_2^f)$. K was defined by

$$K = \frac{\theta^2}{\sigma^2} \sqrt{\int_0^T \int_0^T (\Gamma^{X^0}(s, t))^2 ds dt},$$

where Γ^{X^0} is the covariance function of the Ornstein–Uhlenbeck process starting from 0. Plugging this into equation (A.1) yields

$$\text{Var}\left(\int_0^t f(s) dX_s\right) \geq (1 - \sqrt{K})^2 \text{Var}(G_2^f) = \underbrace{(1 - \sqrt{K})^2 \sigma^2}_{=:c^2} \|f\|_2^2.$$

This is the wanted inequality. \square

A straightforward consequence of Proposition A.1 is that $\|f\|_J = 0 \Leftrightarrow \|f\|_2 = 0$ so that equivalence classes in $\mathcal{E}([0, T])$ for the relation \sim_J are *almost surely* equal functions. Another consequence is that the sets of Cauchy sequences and convergent sequences for the two norms on $E([0, T])$ coincide, and thus the corresponding completions of $E([0, T])$ are the same. In other words, in the case of Ornstein–Uhlenbeck processes that satisfy the condition $\theta T \leq \frac{4}{3}$, we have $F = L^2([0, T])$.

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