# Almost sure approximations in Hölder norms of a general stochastic process defined by a Young integral 

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## Organization of the talk

(1) Introduction and motivation
(2) Approximation of $Y$ paths via wavelets
(3) A better rate of convergence in Wiener chaos framework
(4) Some classes of examples

Let $\{\sigma(s)\}_{s \in \mathbb{R}}$ and $\{X(s)\}_{s \in \mathbb{R}}$ be two real-valued stochastic processes indexed by $\mathbb{R}$ and defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$. We always assume that they satisfy the following assumption:

Fundamental condition: There are two exponents $\alpha, \beta \in(0,1)$ satisfying $\alpha+\beta>1$ such that, for all compact interval $\mathcal{K} \subset \mathbb{R}$ the restrictions of the paths of $\sigma($ resp. $X)$ to $\mathcal{K}$ belong to the Hölder space $C^{\alpha}(\mathcal{K})\left(\right.$ resp. $\left.C^{\beta}(\mathcal{K})\right)$.

This is a very classical condition under which, for each $t \in \mathbb{R}$, the pathwise Young integral

$$
\begin{equation*}
Y(t):=\int_{0}^{t} \sigma(s) \mathrm{d} X(s) \tag{1.1}
\end{equation*}
$$

is a well-defined real-valued random variable. Such kind of process $\{Y(t)\}_{t \in \mathbb{R}}$ is closely connected to (stochastic) differential equations driven by fractional Brownian motions and more generally by (random) Hölder functions with correlated increments (see e.g. Gubinelli, Lejay and Tindel (2007); Gubinelli, Imkeller and Perkowski (2016); Lejay (2010); Lyons, Caruana and Lévy (2004)).
$\Longrightarrow I t$ is useful to find approximation procedures for $Y$ paths which converge at the fastest possible rate.

## A brief recall on the definition of pathwise Young integral

When $t>0$, let
$\left(\mathcal{D}_{n}\right)_{n \in \mathbb{Z}_{+}}:=\left(\left\{\delta_{0}^{n}, \delta_{1}^{n}, \ldots, \delta_{r_{n}}^{n}: 0=\delta_{0}^{n}<\delta_{1}^{n}<\ldots<\delta_{r_{n}}^{n}=t\right\}\right)_{n \in \mathbb{Z}_{+}}$be an arbitrary sequence of partitions of $[0, t]$, such that $\left|\mathcal{D}_{n}\right|:=\max _{1 \leq i \leq r_{n}}\left(\delta_{i}^{n}-\delta_{i-1}^{n}\right) \xrightarrow[n \rightarrow+\infty]{ } 0$. Then, under the previous condition, for all $\omega \in \Omega$, the Riemann-Stieltjes sum:

$$
\begin{equation*}
\sum_{i=1}^{r_{n}} \sigma\left(\delta_{i-1}^{n}, \omega\right)\left(X\left(\delta_{i}^{n}, \omega\right)-X\left(\delta_{i-1}^{n}, \omega\right)\right) \tag{1.2}
\end{equation*}
$$

converges, when $n \rightarrow+\infty$, to a random finite limit not depending on the choice of the sequence of partitions. The pathwise Young integral $\int_{0}^{t} \sigma(s, \omega) \mathrm{d} X(s, \omega)$ is simply defined to be this limit. The case where $t<0$ is similar except that $[0, t]$ has to be replaced by $[t, 0]$ and $\int_{0}^{t}$ means $-\int_{t}^{0}$.
When a compactly supported function $f(\bullet, \omega)$ from $\mathbb{R}$ to $\mathbb{R}$ is $\alpha$-Hölder continuous with a support included in a compact interval $J$, then by convention one sets

$$
\int_{\mathbb{R}} f(s, \omega) \mathrm{d} X(s, \omega):=\int_{J} f(s, \omega) \mathrm{d} X(s, \omega)
$$

Of course, Young integral has much less "good" properties than the usual Lebesgue integral. Yet, it satisfies the following fundamental inequality, whose proof can for instance be found in the well-known book Lyons, Caruana and Lévy (2004) on rough paths.

## Proposition 1.1 (Young-Loeve inequality)

There exists a positive finite constant $\Lambda_{\alpha+\beta}$, depending only on $\alpha+\beta>1$, such that the inequality

$$
\begin{align*}
\mid \int_{t_{1}}^{t_{2}} \sigma(s) \mathrm{d} X(s)-\sigma(\tilde{s})( & \left.X\left(t_{2}\right)-X\left(t_{1}\right)\right) \mid \\
& \leq \Lambda_{\alpha+\beta}\|\sigma\|_{C^{\alpha}\left(\left[t_{1}, t_{2}\right]\right)}\|X\|_{C^{\beta}\left(\left[t_{1}, t_{2}\right]\right)}\left(t_{2}-t_{1}\right)^{\alpha+\beta} \tag{1.3}
\end{align*}
$$

holds for all real numbers $t_{1} \leq \tilde{s} \leq t_{2}$.
$\Longrightarrow$ The restrictions of $Y$ paths to any compact interval $\mathcal{K} \subset \mathbb{R}$ belong to the Hölder space $C^{\beta}(\mathcal{K})$.

## Appoximation of $Y$ paths through Euler scheme

Let $I:=[0, T]$, where $T \in \mathbb{N}$ is fixed once and for all. The classical Euler scheme, corresponding to Riemann-Stieltjes sums associated with dyadic intervals of fixed length $2^{-J}, J \in \mathbb{N}$, provides a natural method for approximating paths of the process $\{Y(t)\}_{t \in l}$, which can be connected to Haar basis. More precisely, for every $m \in\left\{1, \ldots, 2^{J} T\right\}$, one approximates $Y\left(m / 2^{J}\right):=\int_{0}^{m / 2^{J}} \sigma(s) \mathrm{d} X(s)$ by

$$
\begin{equation*}
Y_{J}\left(m / 2^{J}\right):=\int_{0}^{m / 2^{J}}\left(\sum_{l=0}^{m-1} \sigma\left(\tilde{s}_{J, I}\right) \mathbb{1}_{\left[I / 2^{J},(I+1) / 2^{J}\right)}(s)\right) \mathrm{d} X(s)=\sum_{l=0}^{m-1} \sigma\left(\tilde{s}_{J, l}\right) \Delta_{J, I}(X) \tag{1.4}
\end{equation*}
$$

where $\Delta_{J, I}(X):=X\left((I+1) / 2^{J}\right)-X\left(I / 2^{J}\right)$, and $\tilde{s}_{J, I} \in\left[I / 2^{J},(I+1) / 2^{J}\right]$ can be chosen arbitrarily. Thus, taking it such that $\sigma\left(\tilde{s}_{J, I}\right)=2^{J} \int_{I / 2^{J}}^{(I+1) / 2^{J}} \sigma(s) \mathrm{d} s$, one has

$$
\begin{equation*}
Y_{J}\left(m / 2^{J}\right)=\int_{\mathbb{R}} \operatorname{Proj}_{V_{J}^{\mathrm{H}}}\left(\sigma \mathbb{1}_{\left[0, m / 2^{J}\right]}\right)(s) \mathrm{d} X(s), \tag{1.5}
\end{equation*}
$$

where $V_{J}^{\mathrm{H}}:=\overline{\operatorname{span}\left\{\mathbb{1}_{\left[I / 2^{J},(I+1) / 2^{J}\right)}: I \in \mathbb{Z}\right\}}$ is the closed subspace of $L^{2}(\mathbb{R})$ issued from the multiresolution analysis generating the Haar basis.

Once one has the $Y_{J}\left(m / 2^{J}\right)$ 's, using linear interpolation, one gets a random function, from I to $\mathbb{R}, t \mapsto Y_{J}^{R S}(t)$ which approximates the whole path $t \mapsto Y(t)$. More precisely, one sets $Y_{J}^{R S}(0):=0, Y_{J}^{R S}(T):=Y_{J}(T)$, and, for every $t \in i \quad:=(0, T)$,

$$
Y_{J}^{R S}(t):=Y_{J}\left(\frac{\left[2^{J} t\right]}{2^{J}}\right)+\left(2^{J} t-\left[2^{J} t\right]\right)\left(Y_{J}\left(\frac{\left[2^{J} t\right]+1}{2^{J}}\right)-Y_{J}\left(\frac{\left[2^{J} t\right]}{2^{J}}\right)\right) .
$$

The proof of the following proposition manly relies on the Young-Loeve inequality.

## Proposition 1.2

There exists a random finite constant $c>0$ such that for all $\gamma \in[0, \beta)$ and $J \in \mathbb{N}$, one has

$$
\begin{equation*}
\left\|Y-Y_{J}^{R S}\right\|_{C^{\gamma}(I)} \leq c 2^{-J \min (\beta-\gamma, \alpha+\beta-1)} . \tag{1.6}
\end{equation*}
$$

Question: Is it possible to find an approximation procedure for $\{Y(t)\}_{t \in I}$ paths allowing to have a better rate of convergence than the one provided by (1.6)?

Studying this issue is the main motivation of our talk.

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4 Some classes of examples

Our next goal is to provide a generalization of Proposition 1.2 in a wavelet-based framework (see e.g. Meyer (1990, 1992); Daubechies (1992)). The collection of functions, from $\mathbb{R}$ to itself,

$$
\begin{equation*}
\{\varphi(\bullet-I): I \in \mathbb{Z}\} \cup\left\{2^{j / 2} \psi\left(2^{j} \bullet-k\right):(j, k) \in \mathbb{Z}_{+} \times \mathbb{Z}\right\} \tag{2.1}
\end{equation*}
$$

satisfies one of the following two hypotheses.
$\left(\mathcal{H}_{1}\right)$ This collection is simply the Haar basis of $L^{2}(\mathbb{R})$, in other words one has $\varphi:=\mathbb{1}_{[0,1)}$ and $\psi:=\mathbb{1}_{[0,1 / 2)}-\mathbb{1}_{[1 / 2,1)}$.
$\left(\mathcal{H}_{2}\right)$ This collection is an arbitrary compactly supported orthonormal wavelet basis of $L^{2}(\mathbb{R})$ such that the scaling function $\varphi$ and the mother wavelet $\psi$ are $\alpha$-Hölder continuous on $\mathbb{R}$ with a support included in $[-N, N]$, where $N \in \mathbb{N}$. Thus, setting $N_{j, k}^{-}:=(k-N) / 2^{j}$ and $N_{j, k}^{+}:=(k+N) / 2^{j}$, one gets that

$$
\begin{equation*}
\operatorname{supp} \psi\left(2^{j} \bullet-k\right) \subseteq\left[N_{j, k}^{-}, N_{j, k}^{+}\right] . \tag{2.2}
\end{equation*}
$$

It is known that: $\int_{\mathbb{R}} \varphi(x) \mathrm{d} x=1, \int_{\mathbb{R}} \psi(x) \mathrm{d} x=0$, and the integer translates of $\varphi$ form "a partition of unity" i.e. $\sum_{l \in \mathbb{Z}} \varphi(x-l)=1$, for all $x \in \mathbb{R}$.

The increasing sequence $\left(V_{J}\right)_{J \in \mathbb{Z}}$ of closed subspaces of $L^{2}(\mathbb{R})$ denotes the multiresolution analysis associated to this basis, that is

$$
V_{J}:=\overline{\operatorname{span}\left\{\varphi\left(2^{J} \bullet-I\right): I \in \mathbb{Z}\right\}}, \quad \text { for all } J \in \mathbb{Z} .
$$

For all fixed $J \in \mathbb{N}$, let $\left\{Y_{J}^{W}(t)\right\}_{t \in I}$ be the stochastic process defined, for each $t \in I$, as

$$
\begin{align*}
& Y_{J}^{W}(t):=\int_{\mathbb{R}} \operatorname{Proj}_{V_{J}}\left(\sigma \mathbb{1}_{[0, t]}\right)(s) \mathrm{d} X(s)  \tag{2.3}\\
& =\sum_{l \in \mathbb{L}_{J, t} \cup \partial \mathbb{L} J, t}\left(2^{J} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s\right)\left(\int_{N_{J, l}^{-}}^{N_{J, l}^{+}} \varphi\left(2^{J} s-l\right) \mathrm{d} X(s)\right),
\end{align*}
$$

where $\mathbb{L}_{J, t}:=\left\{I \in \mathbb{Z}:\left[N_{j, k}^{-}, N_{j, k}^{+}\right] \subseteq[0, t]\right\}$ and
$\partial \mathbb{L}_{J, t}:=\left\{1 \in \mathbb{Z} \backslash \mathbb{L}_{J, t}:\left[N_{j, k}^{-}, N_{j, k}^{+}\right] \cap[0, t] \neq \emptyset\right\}$. Observe that $\operatorname{card}\left(\mathbb{L}_{J, t}\right) \leq c 2^{J}$
and $\operatorname{card}\left(\partial \mathbb{L}_{\mu, t}\right) \leq c$, where the deterministic constant $c>0$ does not depend on $(J, t)$. The following theorem provides a generalization of Proposition 1.2.

## Theorem 2.1

There is a finite random constant $c>0$ such that, for all $\gamma \in[0, \beta)$ and $J \in \mathbb{N}$, one has

$$
\left\|Y-Y_{J}^{W}\right\|_{C^{\gamma}(I)} \leq c 2^{-J \min (\beta-\gamma, \alpha+\beta-1)}
$$

The proof of Theorem 2.1 mainly relies on the following technical lemma.

## Lemma 2.1

There is a random constant not depending on $(J, I, t)$ such that:

$$
\begin{gather*}
\left|2^{J} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s\right| \leq c, \quad \text { if } I \in \mathbb{Z},  \tag{2.4}\\
\left|2^{J} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s-\sigma\left(I / 2^{j}\right)\right| \leq c 2^{-J \alpha}, \quad \text { if } I \in \mathbb{L}_{J, t},  \tag{2.5}\\
\left|\int_{N_{J, I}^{-}}^{N_{J, I}^{\prime}} \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right| \leq c 2^{-J \beta}, \quad \text { if } I \in \mathbb{Z},  \tag{2.6}\\
\left|\int_{0}^{t}\left(\sigma(s)-\sigma\left(I / 2^{j}\right)\right) \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right| \leq c 2^{-J(\alpha+\beta)}, \quad \text { if } I \in \mathbb{L}_{J, t},
\end{gather*}
$$

and

$$
\begin{equation*}
\left|\int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right| \leq c 2^{-J \beta}, \quad \text { if } I \in \partial \mathbb{L}_{J, t} . \tag{2.8}
\end{equation*}
$$

The main ideas of the proof of Lemma 2.1: The proofs of the first two inequalities are standard and easy. The proofs of the three others are rather similar and rely on the Young-Loeve inequality. So, we only give the one which concerns $\int_{N_{J, I}^{\prime}}^{N_{J, I}^{+}} \varphi\left(2^{\lrcorner} s-I\right) \mathrm{d} X(s)$. Using the Young-Loeve inequality, one gets

$$
\begin{align*}
& \left|\int_{N_{J, l}^{-}}^{N_{J, l}^{+}} \varphi\left(2^{J} s-l\right) \mathrm{d} X(s)-\varphi\left(2^{J} N_{J, l}^{-}-l\right)\left(X\left(N_{J, l}^{+}\right)-X\left(N_{J, l}^{-}\right)\right)\right| \\
& \quad \leq \Lambda_{\alpha+\beta}\left\|\varphi\left(2^{J} \bullet-l\right)\right\|_{C^{\alpha}\left(\left[N_{J, l}^{-}, N_{J, l}^{+}\right)\right.}\|X\|_{C^{\beta}\left(\left[N_{J, l}^{-}, N_{J, l}^{+}\right)\right.}\left(N_{J, l}^{+}-N_{J, l}^{-}\right)^{\alpha+\beta} . \tag{2.9}
\end{align*}
$$

Next, noticing that $\varphi\left(2^{J} N_{J, I}^{-}-I\right)=0\left(\right.$ since $\left.\operatorname{supp} \varphi\left(2^{\lrcorner} \bullet^{-}-I\right) \subseteq\left[N_{J, I}^{-}, N_{J, l}^{+}\right]\right)$,

$$
\left\|\varphi\left(2^{J} \bullet-l\right)\right\|_{C^{\alpha}\left(\left[N_{J, l}^{-}, N_{j, l}, l\right)\right.} \leq 2^{J \alpha}\|\varphi\|_{C^{\alpha}([-N, N])}
$$

and $N_{J, l}^{+}-N_{J, l}^{-}=2^{1-J} N$, one obtains that

$$
\left|\int_{N_{J, I}^{-}}^{N_{J, I}^{+}} \varphi\left(2^{J} s-l\right) \mathrm{d} X(s)\right| \leq c 2^{-J \beta} .
$$

$\square$

Sketch of the proof of Theorem 2.1 in the case $\gamma=0$ : Using the fact that the integer translates of $\varphi$ form "a partition of unity", one gets

$$
\begin{align*}
Y(t) & :=\int_{0}^{t} \sigma(s) \mathrm{d} s=\int_{0}^{t} \sigma(s)\left(\sum_{I \in \mathbb{L}_{J, t} \cup \partial \mathbb{L}, t} \varphi\left(2^{J} s-l\right)\right) \mathrm{d} s \\
& =\sum_{I \in \mathbb{L}_{J, t} \cup \partial \mathbb{L}_{J, t}} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s \\
& \simeq \sum_{I \in \mathbb{L}_{J, t}} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s . \tag{2.10}
\end{align*}
$$

On the other hand, one has

$$
\begin{align*}
& Y_{J}^{W}(t)=\sum_{l \in \mathbb{L}_{, t} \cup \partial \mathbb{L}_{J, t}}\left(2^{J} \int_{0}^{t} \sigma(s) \varphi\left(2^{J} s-I\right) \mathrm{d} s\right)\left(\int_{N_{J, I}^{-}}^{N_{J, I}^{+}} \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right) \\
& \simeq \sum_{I \in \mathbb{L}_{J, t}} \int_{0}^{t} \sigma\left(I / 2^{j}\right) \varphi\left(2^{J} s-I\right) \mathrm{d} X(s) . \tag{2.11}
\end{align*}
$$

One knows from Lemma 2.1 that the approximation errors in (2.10) and (2.11) for the norm $\|\cdot\|_{I, \infty}$ are less than $c_{1} 2^{-J \min (\beta, \alpha+\beta-1)}$. Finally, using the same lemma and the inequality $\operatorname{card}\left(\mathbb{L}_{J, T}\right) \leq c_{2} 2^{J}$, it follows that

$$
\begin{aligned}
\left\|Y-Y_{J}^{W}\right\|_{I, \infty} & \simeq\left\|\sum_{I \in \mathbb{L}_{J, \bullet}} \int_{0}^{\bullet}\left(\sigma(s)-\sigma\left(I / 2^{j}\right)\right) \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right\|_{I, \infty} \\
& \leq \sum_{I \in \mathbb{L}_{J, T}}\left\|\int_{0}^{\bullet}\left(\sigma(s)-\sigma\left(I / 2^{j}\right)\right) \varphi\left(2^{J} s-I\right) \mathrm{d} X(s)\right\|_{I, \infty} \\
& \leq c_{3} \operatorname{card}\left(\mathbb{L}_{J, T}\right) 2^{-J(\alpha+\beta)} \leq c_{4} 2^{-J(\alpha+\beta-1)}
\end{aligned}
$$

One of the main advantages of the wavelet approach is that the difference $Y(t)-Y^{W}(t)$ can be expressed in Hölder spaces in an explicit exploitable way:

$$
\begin{equation*}
Y(t)-Y^{W}(t)=\sum_{j=J}^{+\infty} \sum_{k \in \mathbb{L}_{j, t} \cup \partial \mathbb{L}_{j, t}} a_{j, k}(t) \lambda_{j, k} \tag{2.12}
\end{equation*}
$$

where $a_{j, k}(t):=2^{j} \int_{0}^{t} \sigma(s) \psi\left(2^{j} s-k\right) \mathrm{d} s$ and $\lambda_{j, k}:=\int_{N_{j, k}^{-,}}^{N_{j, k}^{+}} \psi\left(2^{j} s-I\right) \mathrm{d} X(s)$.

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First, it is useful to make some brief recalls on the notion of Wiener chaos; our presentation of it is inspired by the one in the book Janson (1997).

## Definition 3.1 (Wiener chaos)

Let $G$ be an arbitrary fixed Gaussian subspace of $L^{2}(\Omega, \mathcal{F}, \mathbb{P})$, that is a closed subspace consisting of real-valued centred Gaussian random variables.
Let $n \in \mathbb{Z}_{+}$, the Wiener chaos of order $n$ associated with $G$ is denoted by $\overline{\mathcal{P}}_{n}(G)$, or more simply by $\overline{\mathcal{P}}_{n}$.
The space $\overline{\mathcal{P}}_{0}$ is defined to be the closed subspace of $L^{2}(\Omega, \mathcal{F}, \mathbb{P})$ consisting of all the constant random variables.
When $n \geq 1$, the space $\overline{\mathcal{P}}_{n}$ is defined as the closed subspace of $L^{2}(\Omega, \mathcal{F}, \mathbb{P})$ spanned by the following set of random variables:

$$
\left\{\prod_{l=1}^{n} g_{l}^{m_{l}}:\left(g_{1}, \ldots, g_{n}\right) \in G^{n} \text { and }\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}_{+}^{n} \text { with } \sum_{l=1}^{n} m_{l} \leq n\right\}
$$

## Remark 3.1

(a) One clearly has $\overline{\mathcal{P}}_{n} \subseteq \overline{\mathcal{P}}_{n+1}$, for every $n \in \mathbb{Z}_{+}$. Moreover, for all fixed $p \in(0,+\infty)$, the space $\overline{\mathcal{P}}_{*}:=\bigcup_{n \in \mathbb{Z}_{+}} \overline{\mathcal{P}}_{n}$ is dense in $L^{p}(\Omega, \mathcal{F}(G), \mathbb{P})$.
(b) The $L^{p}(\Omega)$-norms are equivalent on $\overline{\mathcal{P}}_{n}$, for all fixed $n \in \mathbb{Z}_{+}$.
(c) For every fixed $n \in \mathbb{Z}_{+}$and for each sequence of random variables in $\overline{\mathcal{P}}_{n}$, convergence in probability is equivalent to convergence in $L^{p}(\Omega)$-norm, for any fixed $p \in(0,+\infty)$.
(d) For all fixed integer $n \geq 1$, there exists a positive finite universal constant $c(n)$, depending only on $n$, such that, for every random variable $\chi \in \overline{\mathcal{P}}_{n}$ and for each real number $y \geq 2$, one has

$$
\begin{equation*}
\mathbb{P}\left(|\chi|>y\|\chi\|_{L^{2}(\Omega)}\right) \leq \exp \left(-c(n) y^{2 / n}\right) \tag{3.13}
\end{equation*}
$$

## Definition 3.2 (the Wiener chaos condition $(\mathcal{W C})$ )

One says that the stochastic process $\{Y(t)\}_{t \in I}$ satisfies $(\mathcal{W C})$ when, for some arbitrary integer $n \geq 1$, the integrand $\{\sigma(s)\}_{s \in \mathbb{R}}$ and the integrator $\{X(s)\}_{s \in \mathbb{R}}$ are two stochastic processes belonging to the Wiener chaos $\overline{\mathcal{P}}_{n}$ (i.e. $\sigma(s) \in \overline{\mathcal{P}}_{n}$ and $X(s) \in \overline{\mathcal{P}}_{n}$, for all $s \in \mathbb{R}$ ) and possessing the following two properties:
$\left(\mathcal{C}_{1}\right)$ There exist $\alpha_{0}, \beta_{0} \in(0,1]$, satisfying $\alpha_{0}+\beta_{0}>1$, such that, on any compact interval $\mathcal{K},\{\sigma(s)\}_{s \in \mathcal{K}}$ and $\{X(s)\}_{s \in \mathcal{K}}$ are respectively $\alpha_{0}$ and $\beta_{0}$ Hölder continuous in the sense of the $L^{2}(\Omega)$-norm.
$\left(\mathcal{C}_{2}\right)$ The "wavelet coefficients" $a_{j, k}:=a_{j, k}(T)$ and $\lambda_{j, k}$ have the following "short-range dependence" property: the inequality

$$
\begin{equation*}
\sum_{k_{1} \in \mathbb{L}_{j, T}} \sum_{k_{2} \in \mathbb{L}_{j, T}}\left|\mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right]\right| \leq c 2^{-2 j\left(\alpha_{0}+\beta_{0}-1 / 2\right)} \tag{3.14}
\end{equation*}
$$

is satisfied, for some finite deterministic constant $c>0$ and for all positive integer $j$ large enough.

## Remark 3.2

(a) The condition $\left(\mathcal{C}_{1}\right)$ implies that paths of $\sigma$ and $X$ respectively belong to the Hölder spaces $C^{\alpha}(\mathcal{K})$ and $C^{\beta}(\mathcal{K})$ for any $\alpha \in\left(0, \alpha_{0}\right)$ and $\beta \in\left(0, \beta_{0}\right)$ (Kolmogorov-Čentsov Hölder continuity theorem).
(b) One has $Y(t) \in \overline{\mathcal{P}}_{2 n}$ for all $t$, since $\sigma(s) \in \overline{\mathcal{P}}_{n}$ and $X(s) \in \overline{\mathcal{P}}_{n}$ for every $s$.
(c) The same argument shows that, for all $(t, j, k)$, the "wavelet coefficients" $a_{j, k}(t)$ and $\lambda_{j, k}$ are in $\overline{\mathcal{P}}_{n}$; therefore $Y_{j}^{W}(t) \in \overline{\mathcal{P}}_{2 n}$.

The main goal of the present section is to obtain the following theorem.

## Theorem 3.1

Under the condition $(\mathcal{W C})$, for any fixed real numbers $\alpha \in\left(0, \alpha_{0}\right), \beta \in\left(0, \beta_{0}\right)$ and $\gamma \geq 0$ satisfying $\alpha+\beta>1$ and $\gamma<\min (\beta, 1 / 2)$, there is a finite random constant $c>0$ such that the following inequality holds almost surely and for each $J \in \mathbb{N}$.

$$
\begin{equation*}
\left\|Y-Y_{J}^{W}\right\|_{C^{\gamma}(I)} \leq c 2^{-J \min (\beta-\gamma, \alpha+\beta-1 / 2-\gamma)} . \tag{3.15}
\end{equation*}
$$

We focus on the case $\gamma=0$, that is:

$$
\begin{equation*}
\left\|Y-Y_{J}^{W}\right\|_{I, \infty} \leq c 2^{-J \min (\beta, \alpha+\beta-1 / 2)} \tag{3.16}
\end{equation*}
$$

In the rest of the section, one presents the main lines of the strategy allowing to derive (3.16).

For each fixed $j \in \mathbb{N}$, one denotes by $Z_{j}:=\left\{Z_{j}(t)\right\}_{t \in I}$ the stochastic process in $\overline{\mathcal{P}}_{2 n}$ with Lipschitz continuous paths defined, for all $t \in I$, as

$$
\begin{equation*}
Z_{j}(t):=\sum_{k \in \mathbb{L}_{j, t} \cup \partial \mathbb{L}_{j, t}} a_{j, k}(t) \lambda_{j, k} \tag{3.17}
\end{equation*}
$$

One already knows from Theorem 2.1 that:

$$
Y(\bullet, \omega)-Y_{j}^{W}(\bullet, \omega)=\sum_{j=J}^{+\infty} Z_{j}(\bullet, \omega)
$$

where the convergence of the series holds in the Hölder space $C^{\gamma}(I)$, for any $\gamma \in[0, \beta)$.

Therefore, using the triangle inequality, one has

$$
\begin{equation*}
\left\|Y(\bullet, \omega)-Y_{j}^{W}(\bullet, \omega)\right\|_{I, \infty} \leq \sum_{j=J}^{+\infty}\left\|Z_{j}(\bullet, \omega)\right\|_{I, \infty} \tag{3.18}
\end{equation*}
$$

Thus, in order to derive $\left\|Y-Y_{J}^{W}\right\|_{I, \infty} \leq c 2^{-J \min (\beta, \alpha+\beta-1 / 2)}$, it is enough to obtain the following lemma.

## Lemma 3.1

One has almost surely

$$
\begin{equation*}
\sup _{j \in \mathbb{N}}\left\{2^{j \min (\beta, \alpha+\beta-1 / 2)}\left\|Z_{j}\right\|_{I, \infty}\right\}<+\infty \tag{3.19}
\end{equation*}
$$

Next, let us point out that $\left\|Z_{j}\right\|_{I, \infty}:=\sup _{t \in I}\left|Z_{j}(t)\right|$ is the supremum of infinitely many random variables. Actually, it is more convenient to work with a supremum of finite number of them; this can be done thanks to the following lemma.

## Lemma 3.2

For each $j \in \mathbb{N}$, one sets

$$
\nu\left(Z_{j}\right):=\sup _{2^{-j} \mid \in I}\left|Z_{j}\left(2^{-j} I\right)\right| .
$$

Then, one has almost surely

$$
\begin{equation*}
\sup _{j \in \mathbb{N}}\left\{2^{j \beta}\left|\left\|Z_{j}\right\|_{I, \infty}-\nu\left(Z_{j}\right)\right|\right\}<+\infty . \tag{3.20}
\end{equation*}
$$

Proof of Lemma 3.2: For all $t \in I$, let $d_{j}(t)$ be the dyadic number of order $j$ such that $t \in\left[d_{j}(t), d_{j}(t)+2^{-j}\right)$.
Then, one has $\mathbb{L}_{j, t}=\mathbb{L}_{j, d_{j}(t)}$ and $a_{j, k}(t)=a_{j, k}\left(d_{j}(t)\right)$ for any $k \in \mathbb{L}_{j, t}$. Therefore

$$
\begin{aligned}
& \left|Z_{j}(t)-Z_{j}\left(d_{j}(t)\right)\right| \leq \sum_{k \in \partial \mathbb{L}_{j, t}}\left|a_{j, k}(t)-a_{j, k}\left(d_{j}(t)\right)\right|\left|\lambda_{j, k}\right| \\
& \leq c_{1} 2^{-j \beta} \sum_{k \in \partial \mathbb{L}_{j, t}} 2^{j} \int_{d_{j}(t)}^{t}\left|\sigma(s) \psi\left(2^{j} s-k\right)\right| \mathrm{d} s \\
& \leq c_{1} 2^{-j \beta}\|\sigma\|_{I, \infty}\|\psi\|_{\infty} \operatorname{card}\left(\partial \mathbb{L}_{j, t}\right) \leq c_{2} 2^{-j \beta}
\end{aligned}
$$

In view of the previous lemma, it turns out that for deriving the theorem it is enough to show that

$$
\begin{equation*}
\sup _{j \in \mathbb{N}}\left\{2^{j \min (\beta, \alpha+\beta-1 / 2)} \nu\left(Z_{j}\right)\right\}<+\infty . \tag{3.21}
\end{equation*}
$$

Notice that if one shows that

$$
\begin{equation*}
\sum_{j=1}^{+\infty} \mathbb{P}\left(2^{j \min (\beta, \alpha+\beta-1 / 2)} \nu\left(Z_{j}\right)>1\right)<+\infty \tag{3.22}
\end{equation*}
$$

then the Borel-Cantelli lemma entails that (3.21) holds.
Using the Markov inequality, one has, for every $j \in \mathbb{N}$,

$$
\begin{equation*}
\mathbb{P}\left(2^{j \min (\beta, \alpha+\beta-1 / 2)} \nu\left(Z_{j}\right)>1\right) \leq 2^{j \min (\beta, \alpha+\beta-1 / 2)} \mathbb{E}\left(\nu\left(Z_{j}\right)\right) . \tag{3.23}
\end{equation*}
$$

## Lemma 3.3

There is a constant $c(n)$, only depending on $n$ the order of the chaos, such that for every $j \in \mathbb{N}$, one has

$$
\begin{equation*}
\mathbb{E}\left(\nu\left(Z_{j}\right)\right) \leq c(n) j^{n / 2} \sup _{2^{-j} I \in I}\left(\mathbb{E}\left(\left|Z_{j}\left(2^{-j} l\right)\right|^{2}\right)\right)^{1 / 2} \tag{3.24}
\end{equation*}
$$

Roughly speaking, the proof of the lemma mainly relies on the fact that one has, for some constant $c^{\prime}$, not depending on ( $j, I$ ), and for all $\tau>0$ large enough

$$
\begin{equation*}
\mathbb{P}\left(\left|Z_{j}\left(2^{-j} l\right)\right|>\tau\right) \leq \exp \left(-c^{\prime} \tau^{2 / n}\right) \tag{3.25}
\end{equation*}
$$

In view of Lemma 3.3 and of the fact that $\alpha \in\left(0, \alpha_{0}\right)$ and $\beta \in\left(0, \beta_{0}\right)$, it turns out that for proving the theorem it is enough that

$$
\begin{equation*}
\sup _{j \in \mathbb{N}} \sup _{2^{-j} \mid \in I}\left\{2^{2 j \min \left(\beta_{0}, \alpha_{0}+\beta_{0}-1 / 2\right)} \mathbb{E}\left(\left|Z_{j}\left(2^{-j} /\right)\right|^{2}\right)\right\}<+\infty . \tag{3.26}
\end{equation*}
$$

## End of the proof of Theorem 3.1

Finally, observe that

$$
\begin{align*}
& \mathbb{E}\left(\left|Z_{j}\left(2^{-j} l\right)\right|^{2}\right)=\mathbb{E}\left(\left|\sum_{k \in \mathbb{L}_{j, t} \cup \partial \mathbb{L}_{j, t}} a_{j, k}(t) \lambda_{j, k}\right|^{2}\right) \\
& \simeq \sum_{k_{1} \in \mathbb{L}_{j, t}} \sum_{k_{2} \in \mathbb{L}_{j, t}} \mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right] \\
& \leq \sum_{k_{1} \in \mathbb{L}_{j, T}} \sum_{k_{2} \in \mathbb{L}_{j, T}}\left|\mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right]\right| . \tag{3.27}
\end{align*}
$$

Using the inequality (3.27) and the condition

$$
\begin{equation*}
\sum_{k_{1} \in \mathbb{L}_{j, T}} \sum_{k_{2} \in \mathbb{L}_{j, T}}\left|\mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right]\right| \leq c 2^{-2 j\left(\alpha_{0}+\beta_{0}-1 / 2\right)} \tag{2}
\end{equation*}
$$

one gets the theorem.

## Organization of the talk

(1) Introduction and motivation
(2) Approximation of $Y$ paths via wavelets
(3) A better rate of convergence in Wiener chaos framework

4 Some classes of examples

This section involves constructing wide classes of examples of real-valued stochastic processes $\sigma$ and $X$ satisfying the Wiener chaos condition ( $\mathcal{W} C$ ).

For the sake of convenience, one assumes that these two processes are independent, centred and given by multiple Itô-Wiener integrals. More precisely, for $\mu=\sigma$ or $\mu=X$, one has, for some $N_{\mu} \in \mathbb{N}$ and for all $s \in \mathbb{R}$,

$$
\begin{equation*}
\mu(s)=\int_{\mathbb{R}^{N_{\mu}}}\left(e^{i s\left(\eta_{1}+\ldots+\eta_{N_{\mu}}\right)}-1\right) g_{\mu}\left(\eta_{1}, \ldots, \eta_{N_{\mu}}\right) \mathrm{d} \widehat{\mathbb{W}}_{\mu}\left(\eta_{1}\right) \ldots \mathrm{d} \widehat{\mathbb{W}}_{\mu}\left(\eta_{N_{\mu}}\right), \tag{4.1}
\end{equation*}
$$

where $d \widehat{\mathbb{W}}_{\mu}$ is the "Fourier transform" of a Brownian measure $d \mathbb{W}_{\mu}$ on $\mathbb{R}$, and $g_{\mu}$ is an arbitrary symmectric complex-valued Borel function such that $\overline{g(\eta)}=g(-\eta)$, for all $\eta \in \mathbb{R}^{N_{\mu}}$, and

$$
\begin{equation*}
\int_{\mathbb{R}^{N_{\mu}}} \min \left(1,\left(\eta_{1}+\ldots+\eta_{N_{\mu}}\right)^{2}\right)\left|g\left(\eta_{1}, \ldots, \eta_{N_{\mu}}\right)\right|^{2} \mathrm{~d} \eta_{1} \ldots \mathrm{~d} \eta_{N_{\mu}}<+\infty . \tag{4.2}
\end{equation*}
$$

These properties of $g_{\mu}$ guarantee the existence of the multiple Itô-Wiener integral in (4.1) and the fact that it is real-valued.

## Remark 4.1

- It is worth mentioning that the well-known Gaussian fractional Brownian motion (fBm) of an arbitrary Hurst parameter $H \in(0,1)$ belongs to this class of processes $\mu$ : in its case one has $N_{\mu}=1$ and

$$
\begin{equation*}
g_{\mu}^{\mathrm{fBm}}(\eta)=c|\eta|^{-H-1 / 2}, \quad \text { for almost all } \eta \in \mathbb{R} \tag{4.3}
\end{equation*}
$$

where $c$ is an arbitrary nonvanishing constant.

- Also we mention that the non-Gaussian Rosenblatt process of an arbitrary parameter $d \in(1 / 4,1 / 2)$ belongs to this same class of processes: in its case one has $N_{\mu}=2$ (second order chaos) and

$$
\begin{equation*}
g_{\mu}^{\operatorname{Ros}}\left(\eta_{1}, \eta_{2}\right)=-i\left(\eta_{1}+\eta_{2}\right)^{-1}\left|\eta_{1} \eta_{2}\right|^{-d}, \quad \text { for almost all }\left(\eta_{1}, \eta_{2}\right) \in \mathbb{R}^{2} \tag{4.4}
\end{equation*}
$$

Next, one denotes by $f_{\mu}$ the even and positive Borel function defined, for each $\xi \in \mathbb{R}$, as

$$
\begin{equation*}
f_{\mu}(\xi):=\left(N_{\mu}\right)!\int_{\mathbb{R}^{N_{\mu}-1}}\left|g_{\mu}\left(\xi-\eta_{2}-\ldots-\eta_{N_{\mu}}, \eta_{2}, \ldots, \eta_{N_{\mu}}\right)\right|^{2} \mathrm{~d} \eta_{2} \ldots \mathrm{~d} \eta_{N_{\mu}} \tag{4.5}
\end{equation*}
$$

with the convention that $f_{\mu}(\xi):=\left|g_{\mu}(\xi)\right|^{2}$, when $N_{\mu}=1$. It can be derived from the properties of $g_{\mu}$, the change of variable $\xi=\eta_{1}+\eta_{2}+\ldots+\eta_{N_{\mu}}$ and the " isometry property" of the multiple Itô-Wiener integral that

$$
\begin{equation*}
\int_{\mathbb{R}} \min \left(1, \xi^{2}\right) f_{\mu}(\xi) d \xi<+\infty \tag{4.6}
\end{equation*}
$$

and, for all $s_{1}, s_{2} \in \mathbb{R}$

$$
\begin{equation*}
\mathbb{E}\left[\mu\left(s_{1}\right) \mu\left(s_{2}\right)\right]=\int_{-\infty}^{+\infty}\left(e^{i s_{1} \xi}-1\right)\left(e^{-i s_{2} \xi}-1\right) f_{\mu}(\xi) \mathrm{d} \xi \tag{4.7}
\end{equation*}
$$

Thus the function $f_{\mu}$ can be viewed as a spectral density.

Notice that (4.7) is equivalent to

$$
\begin{equation*}
\mathbb{E}\left[\left|\mu\left(s_{1}\right)-\mu\left(s_{2}\right)\right|^{2}\right]=\mathbb{E}\left[\left|\mu\left(\left|s_{1}-s_{2}\right|\right)\right|^{2}\right]=4 \int_{-\infty}^{+\infty} \sin ^{2}\left(\frac{\left|s_{1}-s_{2}\right| \xi}{2}\right) f_{\mu}(\xi) \mathrm{d} \xi . \tag{4.8}
\end{equation*}
$$

## Remark 4.2

Assume that the process $\{\mu(s)\}_{s \in \mathbb{R}}$ is self-similar of order $\gamma_{0} \in(0,1)$, that is the processes $\{\mu(a s)\}_{s \in \mathbb{R}}$ and $\left\{a^{\gamma_{0}} \mu(s)\right\}_{s \in \mathbb{R}}$ have the same finite-dimensional distributions, for any fixed positive real number a. Then, the corresponding spectral density $f_{\mu}$ satisfies

$$
\begin{equation*}
f_{\mu}(\xi)=c|\xi|^{-2 \gamma_{0}-1}, \quad \text { for almost all } \xi \in \mathbb{R} \tag{4.9}
\end{equation*}
$$

where $c$ is some positive constant. We recall in passing that the Gaussian fractional Brownian motion of Hurst parameter $H \in(0,1)$ is self-similar of order $\gamma_{0}=H$. Also, we recall that the non-Gaussian Rosenblatt process of parameter $d \in(1 / 4,1 / 2)$ is self-similar of order $\gamma_{0}=2 d$.

## Remark 4.3

A sufficient condition for the processes $\sigma$ and $X$ to satisfy $\left(\mathcal{C}_{1}\right)$ i.e. to be, on any compact interval $\mathcal{K}, \alpha_{0}$ and $\beta_{0}$ Hölder continuous for the norm $\|\cdot\|_{L^{2}(\Omega)}$ is the following: there exist two positive finite deterministic constants $c$ and $\xi_{0}$, such that the inequalities

$$
\begin{equation*}
f_{\sigma}(\xi) \leq c|\xi|^{-2 \alpha_{0}-1} \quad \text { and } \quad f_{X}(\xi) \leq c|\xi|^{-2 \beta_{0}-1} \tag{4.10}
\end{equation*}
$$

hold for almost all real number $\xi$ satisfying $|\xi| \geq \xi_{0}$.

## Remark 4.4

Suppose there are a finite constant $c>0$ and two nonnegative integers $U_{0}$ and $V_{0}$ satisfying $U_{0}+V_{0}=2$, such that, for every $j \in \mathbb{N}$ and for all $k_{1}, k_{2} \in \mathbb{L}_{j, T}$, one has

$$
\begin{equation*}
\left|\mathbb{E}\left[a_{j, k_{1}} a_{j, k_{2}}\right]\right| \leq c 2^{-2 j \alpha_{0}}\left(1+\left|k_{1}-k_{2}\right|\right)^{-U_{0}} \tag{4.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\mathbb{E}\left[\lambda_{j, k_{1}} \lambda_{j, k_{2}}\right]\right| \leq c 2^{-2 j \beta_{0}}\left(1+\left|k_{1}-k_{2}\right|\right)^{-V_{0}} . \tag{4.12}
\end{equation*}
$$

Then (C2) is satisfied.

Proof of Remark 4.4: The fact that $\sigma$ and $X$ are independent implies that the associated sequences of "wavelet coefficients" $\left(a_{j, k}\right)_{j, k}$ and $\left(\lambda_{j, k}\right)_{j, k}$ are independent as well. This together with the inequalities (4.11), (4.12) and the fact that $U_{0}+V_{0}=2$ yields

$$
\begin{align*}
& \left|\mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right]\right|=\left|\mathbb{E}\left[a_{j, k_{1}} a_{j, k_{2}}\right]\right|\left|\mathbb{E}\left[\lambda_{j, k_{1}} \lambda_{j, k_{2}}\right]\right| \\
& \leq\left(c 2^{-2 j \alpha_{0}}\left(1+\left|k_{1}-k_{2}\right|\right)^{-U_{0}}\right)\left(c 2^{-2 j \beta_{0}}\left(1+\left|k_{1}-k_{2}\right|\right)^{-V_{0}}\right) \\
& \leq c^{2} 2^{-2 j\left(\alpha_{0}+\beta_{0}\right)}\left(1+\left|k_{1}-k_{2}\right|\right)^{-2}, \tag{4.13}
\end{align*}
$$

for all $k_{1}, k_{2} \in \mathbb{L}_{j, T}$. Then setting $c_{1}:=2 c^{2} \sum_{q \in \mathbb{N}} q^{-2}$ and using (4.13) and the inequality $\operatorname{card}\left(\mathbb{L}_{j, T}\right) \leq c_{2} 2^{j}$, one gets:

$$
\begin{aligned}
& \sum_{k_{1} \in \mathbb{L}_{j, T}} \sum_{k_{2} \in \mathbb{L}_{j, T}}\left|\mathbb{E}\left[a_{j, k_{1}} \lambda_{j, k_{1}} a_{j, k_{2}} \lambda_{j, k_{2}}\right]\right| \\
& \leq c^{2} 2^{-2 j\left(\alpha_{0}+\beta_{0}\right)} \sum_{k_{1} \in \mathbb{L}_{j, T}} \sum_{k_{2} \in \mathbb{L}_{j, T}}\left(1+\left|k_{1}-k_{2}\right|\right)^{-2} \\
& \leq c_{1} \operatorname{card}\left(\mathbb{L}_{j, T}\right) 2^{-2 j\left(\alpha_{0}+\beta_{0}\right)} \leq c_{1} c_{2} 2^{-2 j\left(\alpha_{0}+\beta_{0}-1 / 2\right)}
\end{aligned}
$$

## Proposition 4.1

Assume that the wavelet $\psi$ is the Haar function, that is $\psi:=\mathbb{1}_{[0,1 / 2)}-\mathbb{1}_{[1 / 2,1)}$. Then (4.12) holds as soon as $f_{X}$ is $V_{0}$ times continuously differentiable on $\mathbb{R} \backslash\{0\}$ and satisfies the following condition: There exist two finite deterministic constants $\beta_{0}^{\prime} \in\left[\beta_{0}, 1\right)$ and $c>0$ such that, for all $n \in\left\{0, \ldots, V_{0}\right\}$ and $\xi \in \mathbb{R} \backslash\{0\}$, one has

$$
\begin{equation*}
\left|f_{X}^{(n)}(\xi)\right| \leq c \max \left(|\xi|^{-2 \beta_{0}-n-1},|\xi|^{-2 \beta_{0}^{\prime}-n-1}\right) \tag{1,X}
\end{equation*}
$$

## Proposition 4.2

Let $M \in \mathbb{N}$ be arbitrary and fixed. Assume that the wavelet $\psi$ is continuously differentiable on the real line and has at least $M$ vanishing moments. Then (4.12) holds as soon as $f_{X}$ is $V_{0}$ times continuously differentiable on $\mathbb{R} \backslash\{0\}$ and satisfies the following condition $\left(\mathcal{D}_{M, X}\right)$, which is weaker than $\left(\mathcal{D}_{1, X}\right)$ : There exist two finite deterministic constants $\beta_{0}^{\prime} \in\left[\beta_{0}, 1\right)$ and $c>0$ such that, for all $n \in\left\{0, \ldots, V_{0}\right\}$ and $\xi \in \mathbb{R} \backslash\{0\}$, one has

$$
\left|f_{X}^{(n)}(\xi)\right| \leq c \max \left(|\xi|^{-2 \beta_{0}-n-1},|\xi|^{-2 \beta_{0}^{\prime}-n M-1}\right)
$$

## Proposition 4.3

Assume that the wavelet $\psi$ is the Haar function. Also assume that the integer $U_{0}$ in (4.11) belongs to the set $\{0,1\}$. Then (4.11) holds as soon as $f_{\sigma}$ is $U_{0}$ times continuously differentiable on $\mathbb{R} \backslash\{0\}$ and satisfies the following condition: There exist two finite deterministic constants $\alpha_{0}^{\prime} \in\left[\alpha_{0}, 1\right)$ and $c>0$ such that, for all $n \in\left\{0, \ldots, U_{0}\right\}$ and $\xi \in \mathbb{R} \backslash\{0\}$, one has

$$
\left|f_{\sigma}^{(n)}(\xi)\right| \leq c \max \left(|\xi|^{-2 \alpha_{0}-n-1},|\xi|^{-2 \alpha_{0}^{\prime}-n-1}\right)
$$

## Proposition 4.4

Let $M \in \mathbb{N}$ be arbitrary and fixed. Assume that the wavelet $\psi$ is continuously differentiable on the real line and has at least $M+1$ vanishing moments. Then (4.11) holds as soon as $f_{\sigma}$ is $U_{0}$ times continuously differentiable on $\mathbb{R} \backslash\{0\}$ and satisfies the following condition ( $\mathcal{D}_{M, \sigma}$ ), which is weaker than $\left(\mathcal{D}_{1, \sigma}\right)$ : There exist two finite deterministic constants $\alpha_{0}^{\prime} \in\left[\alpha_{0}, 1\right)$ and $c>0$ such that, for all $n \in\left\{0, \ldots, U_{0}\right\}$ and $\xi \in \mathbb{R} \backslash\{0\}$, one has

$$
\left|f_{\sigma}^{(n)}(\xi)\right| \leq c \max \left(|\xi|^{-2 \alpha_{0}-n-1},|\xi|^{-2 \alpha_{0}^{\prime}-n M-1}\right)
$$

## Remark 4.5

For $\mu=\sigma$ or $\mu=X$, it is clear that $\left(\mathcal{D}_{1, \mu}\right)$ holds when the process $\mu$ is self-similar of order $\gamma_{0}=\alpha_{0}$ or $\gamma_{0}=\beta_{0}$.

## Remark 4.6

A major motivation for weakening the condition ( $\mathcal{D}_{1, \mu}$ ) to the condition ( $\mathcal{D}_{M, \mu}$ ) is the following: the behavior of $f_{\mu}$ in the neighborhood of 0 can then be much more singular, namely $f_{\mu}$ can have infinitely many oscillations in the vicinity of 0 . This is for instance the case, when $f_{\mu}$ is the "chirp function": for all $\xi \in \mathbb{R} \backslash\{0\}$,

$$
\begin{equation*}
f_{\mu}(\xi)=|\xi|^{-2 u-1}+|\xi|^{-2 v-1} \sin ^{2}\left(|\xi|^{-w}\right), \tag{4.14}
\end{equation*}
$$

where the three parameters $u, v$ and $w$ are arbitrary real numbers such that $0<u \leq v<1$ and $w>0$. Observe that the larger is $w$ the more oscillating is this function $f_{\mu}$ in the neighborhood of 0 . Also observe that this function fails to satisfy ( $\mathcal{D}_{1, \mu}$ ); yet, for any integer $M \geq 1+w$, it satisfies ( $\mathcal{D}_{M, \mu}$ ), with $\beta_{0}=u$ and $\beta_{0}^{\prime}=v$.

## Conclusion

## Remark 4.7

Assume for instance that the integrand $\sigma$ and the integrator $X$ are two independent Gaussian fractional Brownian motions whose Hurst parameters satisfy $H_{1} \geq 1 / 2$ and $H_{1}+H_{2}>1$. Then, one has $\alpha_{0}=H_{1}, \beta_{0}=H_{2}$ and $H_{2}=\min \left(\beta_{0}, \alpha_{0}+\beta_{0}-1 / 2\right)$. Thus, it results from Theorem 3.1 that, for all fixed $\gamma \in\left[0, \min \left(H_{2}, 1 / 2\right)\right)$ and for every $\varepsilon>0$ small enough, one has almost surely

$$
\limsup _{J \rightarrow+\infty} 2^{J\left(H_{2}-\gamma-\varepsilon\right)}\left\|Y-Y_{J}^{W}\right\|_{C^{\gamma}(I)}<+\infty
$$

## Proposition 4.5

The rate of convergence provided by (4.15) is optimal, namely: for all fixed $\gamma \in\left[0, \min \left(H_{2}, 1 / 2\right)\right)$ and for every $\varepsilon>0$ small enough, one has almost surely

$$
\begin{equation*}
\limsup _{J \rightarrow+\infty} 2^{J\left(H_{2}-\gamma+\varepsilon\right)}\left\|Y-Y_{J}^{W}\right\|_{C^{\gamma}(I)}=+\infty . \tag{4.16}
\end{equation*}
$$

Proposition 4.5 can be extended to a general framework.

