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# Convergence rate of strong approximations of compound random maps 

Emmanuel Gobet* Mohamed Mrad ${ }^{\dagger}$

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#### Abstract

We consider a random map $x \mapsto F(\omega, x)$ and a random variable $\Theta(\omega)$, and we denote by $F^{N}(\omega, x)$ and $\Theta^{N}(\omega)$ their approximations: We establish a strong convergence result, in $\mathbf{L}_{p}$-norms, of the compound approximation $F^{N}\left(\omega, \Theta^{N}(\omega)\right)$ to the compound variable $F(\omega, \Theta(\omega)$ ), in terms of the approximations of $F$ and $\Theta$. Two applications of this result are then developed: Firstly, composition of two Stochastic Differential Equations through their initial conditions; secondly, approximation of stochastic processes (possibly non semi-martingales) at random times (possibly non stopping times).


Keywords: strong approximation, Garsia-Rodemich-Rumsey lemma, Euler schemes, Iterated Brownian motion, local time, Fractional Brownian motion

MSC: $60 \mathrm{Hxx}, 60 \mathrm{Gxx}$

## 1 Introduction

Since the seventies, the numerical analysis of stochastic systems is a research field on its own and it has tremendous applications in engineering sciences. This work enriches this vast area by addressing the following natural questions. Consider a continuous random map $x \mapsto F(\omega, x)$ and a random variable $\Theta(\omega)$, and their numerical approximations $F^{N}(\omega, x)$ and $\Theta^{N}(\omega)$ for some convergence parameter $N \rightarrow+\infty$ :

- Under which assumptions does the compound approximation $\omega \mapsto F^{N}\left(\omega, \Theta^{N}(\omega)\right)$ converge in $\mathbf{L}_{p}$ to the compound map $\omega \mapsto F(\omega, \Theta(\omega))$ ?
- What is the convergence rate and how does it depend on those related to the approximations $F^{N}$ to $F$ and $\Theta^{N}$ to $\Theta$ ?

[^0]It is easy to guess that the analysis would be straightforward if $\left(F, F^{N}\right)$ were independent of $\left(\Theta, \Theta^{N}\right)$, by using a conditioning argument. On the contrary, here our aim is to allow arbitrary dependencies and study the strong convergence in this general setting (convergence in $\mathbf{L}_{p}$-norms).

Among the applied probability community, there is an increasing interest for strong convergence rates because they constitute the corner stone for designing efficient Multi-Level Monte Carlo methods [Hei01, Gil08] (which significantly speeds-up Crude Monte Carlo methods). In this work, we provide generic results which pave the way for establishing strong convergence rates in complicated situations where results were not available so far. Hopefully, it will open the door for many other interesting issues.

The paper is organized as follows. In Section 2 we state a general convergence result (Theorem 1) estimating the $\mathbf{L}_{p}$-error $\left\|F^{N}\left(\Theta^{N}\right)-F(\Theta)\right\|_{\mathbf{L}_{p}}$, and then we prove it. For this we assume locally uniform approximations on $F^{N}-F$, and local-Hölder continuity on $F$. These assumptions being possibly difficult to check in practice, we then give much easier conditions that imply the first ones, using the Garsia-Rodemich-Rumsey lemma with precise quantitative controls.
In Section 3, we study the error induced by compound Euler schemes related to Stochastic Differential Equations (SDEs for short), through their initial conditions. This question originates in the resolution of Stochastic PDEs using stochastic flows. In Section 4, we analyse the error arising when stochastic processes are evaluated at random times, both being approximated. Then, examples are developed, such as Brownian local times at random points, Fractional Brownian motions or diffusion processes at Brownian time.

## $2 \quad \mathbf{L}_{p}$-approximation of compound random maps

The section is devoted to stating and proving a general result (Theorem 1). Applications are postponed to subsequent sections.

### 2.1 Assumptions

Let $(\mathcal{E},||$.$) be a separable Banach space and (\Omega, \mathcal{F}, \mathbb{P})$ be a probability space. We are given

- a random field, i.e. a $\mathcal{F} \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$-measurable mapping $(\omega, x) \in\left(\Omega, \mathbb{R}^{d}\right) \mapsto$ $F(\omega, x) \in \mathcal{E}$, continuous in $x$ for a.e. $\omega$;
- a $\mathcal{F}$-random variable $\Theta: \Omega \mapsto \mathbb{R}^{d}$.

Let $F^{N}$ and $\Theta^{N}$ be respectively approximations of $F$ and $\Theta$, where $N \rightarrow+\infty$ is a asymptotic parameter; we aim at controlling in $\mathbf{L}_{p}$ the random variable

$$
\omega \in \Omega \mapsto F^{N}\left(\omega, \Theta^{N}(\omega)\right)-F(\omega, \Theta(\omega)) \in \mathcal{E}
$$

which will be denoted by $F^{N}\left(\Theta^{N}\right)-F(\Theta)$ for the sake of simplicity. For $p>0$ and for a random variable $Z: \Omega \mapsto \mathcal{E}$ or $\mathbb{R}^{d}$, we set $\|Z\|_{\mathbf{L}_{p}}=\left(\mathbb{E}|Z|^{p}\right)^{1 / p}$ : We say that
$Z \in \mathbf{L}_{p}$ if $\|Z\|_{\mathbf{L}_{p}}<+\infty$. Despite $\|\cdot\|_{\mathbf{L}_{p}}$ is not a norm for $p<1$, we refer to it as $\mathbf{L}_{p}$-norm to simplify the discussion.
(H1) For any $p>0$, there exist constants $\alpha_{p}^{(\mathbf{H} 1)} \in[0,+\infty)$ and $C_{p}^{(\mathbf{H 1} 1)} \in[0,+\infty)$ such that

$$
\begin{equation*}
\left\|\sup _{|x| \leq \lambda}|F(\cdot, x)|\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(\mathbf{H} 1)} \lambda^{\alpha_{p}^{(\mathbf{H} 1)}}, \quad \forall \lambda \geq 1 \tag{H1}
\end{equation*}
$$

(H2) There is a $\kappa \in(0,1]$ such that for any $p>0$, there exist constants $\alpha_{p}^{(\mathbf{H 2})} \in$ $[0,+\infty)$ and $C_{p}^{(\mathbf{H 2})} \in[0,+\infty)$ such that

$$
\begin{equation*}
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{|F(\cdot, y)-F(\cdot, x)|}{|y-x|^{\kappa}}\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(\mathbf{H} 2)} \lambda^{\alpha_{p}^{(\mathbf{H} 2)}}, \quad \forall \lambda \geq 1 \tag{H2}
\end{equation*}
$$

(H3) For any $p>0$, there exist a constant $\alpha_{p}^{(\mathbf{H 3})} \in[0,+\infty)$ and a sequence $\left(C_{p}^{N,(\mathbf{H 3})}\right)_{N \geq 1}$ with $C_{p}^{N,(\mathbf{H} 3)} \in[0,+\infty)$ such that

$$
\begin{equation*}
\left\|\sup _{|x| \leq \lambda}\left|F^{N}(\cdot, x)-F(\cdot, x)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p}^{N,(\mathbf{H} \mathbf{3})} \lambda^{\alpha_{p}^{(\mathbf{H} \mathbf{3})}}, \quad \forall \lambda \geq 1, \forall N \geq 1 \tag{H3}
\end{equation*}
$$

(H4) For any $p>0$, there exist a constant $C_{p}^{(\mathbf{H} 4-\mathrm{a})} \in[0,+\infty)$ and a sequence $\left(C_{p}^{N,(\mathbf{H} 4-\mathrm{b})}\right)_{N \geq 1}$ with $C_{p}^{N,(\mathbf{H} 4-\mathrm{b})} \in[0,+\infty)$ such that

$$
\begin{align*}
\|\Theta\|_{\mathbf{L}_{p}} \vee\left\|\Theta^{N}\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(\mathbf{H} 4-\mathrm{a})}, \quad \forall N \geq 1  \tag{H4-a}\\
\left\|\Theta^{N}-\Theta\right\|_{\mathbf{L}_{p}} \leq C_{p}^{N,(\mathbf{H} 4-\mathrm{b})}, \quad \forall N \geq 1 \tag{H4-b}
\end{align*}
$$

These conditions state that all random variables belong to any $\mathbf{L}_{p}$, with some locally uniform estimates w.r.t. the space dependance; the extension to belonging to some $\mathbf{L}_{p}$ only would be easy and is left to the reader.

### 2.2 Main results

Had the random variable $\Theta$ be bounded by a finite constant $\Lambda$, we would have directly obtained $\left\|F^{N}(\Theta)-F(\Theta)\right\|_{\mathbf{L}_{p}} \leq C_{p}^{N,(\mathbf{H} 3)} \Lambda^{\alpha_{p}^{(\mathbf{H} 3)}}$. The extension to non bounded r.v. $\Theta$ is non trivial and is being achieved in Theorem 1 and its proof. The following result (inspired by [KS97, Lemma 2.1]) is instrumental in our analysis. In particular, it enables to justify that the quantities of study are well defined as $\mathbf{L}_{p}$ random variables.

Proposition 1. Let $E$ be an Euclidean space. Let $G$ be a $\mathcal{F} \otimes \mathcal{B}(E)$-measurable mapping taking values in $\mathcal{E}$ such that, for any $p>0$ there exist constants $\alpha_{p}^{(G)} \in$ $[0,+\infty)$ and $C_{p}^{(G)} \in[0,+\infty)$ for which

$$
\begin{equation*}
\left\|\sup _{|x| \leq \lambda}|G(\cdot, x)|\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(G)} \lambda^{\alpha_{p}^{(G)}}, \quad \forall \lambda \geq 1 \tag{1}
\end{equation*}
$$

Let $\xi$ be a random variable taking values in $E$, with finite $\mathbf{L}_{p}$ norms for any $p>0$. Then for any $p>0, \omega \mapsto G(\omega, \xi(\omega)) \in \mathbf{L}_{p}$ and for any finite conjugate exponents $r$ and $s\left(r^{-1}+s^{-1}=1\right)$, we have the estimate

$$
\|G(\xi)\|_{\mathbf{L}_{p}} \leq C_{p r}^{(G)}(\zeta(r))^{1 /(p r)} 2^{\alpha_{p r}^{(G)}+1 / p}\left(1+\|\xi\|_{\mathbf{L}_{s\left(\alpha_{p r}\right.}^{(G)} p_{p+1)}}^{\alpha_{p}^{(G)}+1 / p}\right)
$$

where $\zeta(r):=\sum_{n \geq 1} n^{-r}$ is the Riemann zeta function.
The above result will be extended later in Proposition 8 (Subsection 4.4), when the polynomial growth (1) is replaced by an exponential one and when the random variable $\xi$ has exponential moments.

Proof. Using twice Hölder inequalities, we obtain

$$
\begin{aligned}
\mathbb{E}\left(|G(\cdot, \xi)|^{p}\right) & \leq \sum_{n \geq 1} \mathbb{E}\left(\sup _{|x| \leq n}|G(\cdot, x)|^{p} 1_{n-1 \leq|\xi|<n}\right) \\
& \leq \sum_{n \geq 1}\left(\mathbb{E}\left(\sup _{|x| \leq n}|G(\cdot, x)|^{p r}\right)\right)^{1 / r} \mathbb{P}(n-1 \leq|\xi|<n)^{1 / s} \\
& \leq\left[C_{p r}^{(G)}\right]^{p} \sum_{n \geq 1} \frac{1}{n} n^{\alpha_{p r}^{(G)} p+1} \mathbb{P}(n-1 \leq|\xi|<n)^{1 / s} \\
& \leq\left[C_{p r}^{(G)}\right]^{p}\left(\sum_{n \geq 1} \frac{1}{n^{r}}\right)^{1 / r}\left(\sum_{n \geq 1} n^{s\left(\alpha_{p r}^{(G)} p+1\right)} \mathbb{P}(n \leq|\xi|+1<n+1)\right)^{1 / s} \\
& \leq\left[C_{p r}^{(G)}\right]^{p}(\zeta(r))^{1 / r}\left(\mathbb{E}\left((|\xi|+1)^{s\left(\alpha_{p r}^{(G)} p+1\right)}\right)\right)^{1 / s}
\end{aligned}
$$

Therefore, $\|G(\xi)\|_{\mathbf{L}_{p}} \leq C_{p r}^{(G)}(\zeta(r))^{1 /(p r)}\left(1+\|\xi\|_{\mathbf{L}_{s\left(\alpha_{p r}^{(G)}{ }_{p+1)}\right.}}\right)^{\alpha_{p r}^{(G)}+1 / p} \quad$ where we have used the Minkowsky inequality. We complete our statement by using

$$
\begin{equation*}
(a+b)^{\gamma} \leq 2^{(\gamma-1)+}\left(a^{\gamma}+b^{\gamma}\right) \leq 2^{\gamma}\left(a^{\gamma}+b^{\gamma}\right) \tag{2}
\end{equation*}
$$

for any non-negative $a, b, \gamma$.
As a direct consequence of the above result, we deduce that $F(\Theta)$ is any $\mathbf{L}_{p}$ (owing to ( $\mathbf{H} \mathbf{1}$ ) and ( $\mathbf{H} 4-\mathrm{a})$ ). Moreover we can also apply it to $G=F^{N}$ and $\xi=\Theta^{N}$ in view of (H4-a) and since (1) is satisfied (owing to (H1) and (H3)): Thus, $F^{N}\left(\Theta^{N}\right)$ also belongs to any $\mathbf{L}_{p}$.

Our main result below states an error estimate on the approximation of $F(\Theta)$ by $F^{N}\left(\Theta^{N}\right)$, as a function of $N$, through the sequences $\left(C^{N,(\mathbf{H} 3)}\right)_{N \geq 1}$ and $\left(C_{.}^{N,(\mathbf{H} 4-\mathrm{b})}\right)_{N \geq 1}$.

Theorem 1. Assume (H1)-(H2)-(H3)-(H4-a)-(H4-b). Then for any $p>0$ and any $p_{2}>p$, there is a constant $c_{(3)}$ independent on $N$ such that

$$
\begin{equation*}
\left\|F^{N}\left(\Theta^{N}\right)-F(\Theta)\right\|_{\mathbf{L}_{p}} \leq c_{(3)}\left(C_{2 p}^{N,(\mathbf{H 3})}+\left[C_{\kappa p_{2}}^{N,(\mathbf{H} 4-\mathrm{b})}\right]^{\kappa}\right), \quad \forall N \geq 1 \tag{3}
\end{equation*}
$$

Quite intuitively, the global approximation error inherits from that on $F$ and that on $\Theta$ modified by the local Hölder regularity of $x \mapsto F(\omega, x)$.

Proof. Write $F^{N}\left(\Theta^{N}\right)-F(\Theta)=\left[F^{N}\left(\Theta^{N}\right)-F\left(\Theta^{N}\right)\right]+\left[F\left(\Theta^{N}\right)-F(\Theta)\right]$. First, a direct application of Proposition 1 (for $r=s=2$ ) with (H3) and (H4-a) yields $\left\|F^{N}\left(\Theta^{N}\right)-F\left(\Theta^{N}\right)\right\|_{\mathbf{L}_{p}} \leq C_{2 p}^{N,(\mathbf{H} \mathbf{3})}(\zeta(2))^{1 /(2 p)} 2^{\alpha_{2 p}^{(\mathbf{H} \mathbf{3})}+1 / p}\left(1+\left\|\Theta^{N}\right\|_{\mathbf{L}_{2\left(\alpha_{2 p}\right.}^{(\mathbf{H} 3)}{ }_{p+1)}}^{\alpha_{2}^{(\mathbf{H} \mathbf{3})}+1 / p}\right)$

$$
\leq C_{2 p}^{N,(\mathbf{H} \mathbf{3})}(\zeta(2))^{1 /(2 p)} 2^{\alpha_{2 p}^{(\mathbf{H} 3)}+1 / p}\left(1+\left[C_{2\left(\alpha_{2 p}^{(\mathbf{H} 3)} p+1\right)}^{(\mathbf{H} 4-\mathrm{a})}\right]^{\alpha_{2 p}^{(\mathbf{H} 3)}+1 / p}\right)
$$

Consider now the second term $F\left(\Theta^{N}\right)-F(\Theta)$ : Set

$$
H_{\kappa}(\omega, \lambda):=\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{|F(\omega, y)-F(\omega, x)|}{|y-x|^{\kappa}}
$$

and write $\left|F\left(\Theta^{N}\right)-F(\Theta)\right| \leq H_{\kappa}\left(\left|\Theta^{N}\right| \vee|\Theta|\right)\left|\Theta^{N}-\Theta\right|^{\kappa}$. Then the Hölder inequality with $p$-conjugate numbers $\left(p_{1}, p_{2}\right)$ (i.e. $p_{1}^{-1}+p_{2}^{-1}=p^{-1}$ ) gives

$$
\left\|F\left(\Theta^{N}\right)-F(\Theta)\right\|_{\mathbf{L}_{p}} \leq\left\|H_{\kappa}\left(\left|\Theta^{N}\right| \vee|\Theta|\right)\right\|_{\mathbf{L}_{p_{1}}}\left\|\Theta^{N}-\Theta\right\|_{\mathbf{L}_{\kappa p_{2}}}^{\kappa}
$$

The first factor is upper bound using Proposition 1 (for $r=s=2$ ) with (H2) and (H4-b), it readily leads to

$$
\begin{aligned}
& \left\|F\left(\Theta^{N}\right)-F(\Theta)\right\|_{\mathbf{L}_{p}} \\
& \quad \leq\left\|H_{\kappa}\left(\left|\Theta^{N}\right| \vee|\Theta|\right)\right\|_{\mathbf{L}_{p_{1}}}\left\|\Theta^{N}-\Theta\right\|_{\mathbf{L}_{\kappa p_{2}}}^{\kappa} \\
& \quad \leq C_{2 p_{1}}^{(\mathbf{H} 2)}(\zeta(2))^{1 /\left(2 p_{1}\right)} 2^{\alpha_{2 p_{1}}^{(\mathbf{H} 2)}+1 / p_{1}}\left(1+\left\|\left|\Theta^{N}\right| \vee|\Theta|\right\|_{\left.\mathbf{L}_{2\left(\alpha_{2}\right.}^{(\mathbf{H} 2)} p_{1}+1\right)}^{\alpha_{2}^{(\mathbf{H} 2)}+1 / p_{1}}\right)\left[C_{\kappa p_{2}}^{N,(\mathbf{H} 4-\mathrm{b})}\right]^{\kappa} \\
& \quad \leq C_{2 p_{1}}^{(\mathbf{H} 2)}(\zeta(2))^{1 /\left(2 p_{1}\right)} 2^{\alpha_{2 p_{1}}^{(\mathbf{H} 2)}+1 / p_{1}}\left(1+\left[2 C_{2\left(\alpha_{2 p_{1}}^{(\mathbf{H} 2)} p_{1}+1\right)^{(\mathbf{H} 4-\mathrm{a})}}^{\alpha_{2 p_{1}}^{(\mathbf{H 2 )}}+1 / p_{1}}\right)\left[C_{\kappa p_{2}}^{N,(\mathbf{H} 4-\mathrm{b})}\right]^{\kappa} .\right.
\end{aligned}
$$

We are done.

### 2.3 Simplified assumptions

In some situations, checking the assumptions (H1-H2-H3) may be difficult since we evaluate the $\mathbf{L}_{p}$-norms of a maximum. When $x$ is a time variable, we may rely on Doob inequalities and other martingale estimates to achieve this. In other situations, it becomes much more complicated. One can apply the general Kolmogorov continuity criterion for random fields [Kun97, Theorem 1.4.1 p.31], but it does not yield the quantitative estimates we are looking for, in particular regarding the polynomial growth factor in ( $\mathbf{H} \mathbf{1} \mathbf{- H 2} \mathbf{- H 3}$ ). Alternatively, here we use the Garsia-RodemichRumsey lemma [GRR70] (see for instance [Nua06, p.353-354]) which gives refinement compared to the Kolmogorov criterion. This approach has been extensively developed in [BY82] for studying regularity of local times of continuous martingales w.r.t. the space variable.

Lemma 1 (Garsia-Rodemich-Rumsey lemma, control of modulus of continuity). Let $\rho, \Psi: \mathbb{R}_{+} \longrightarrow \mathbb{R}_{+}$be continuous and strictly increasing functions vanishing at zero and such that $\lim _{t \rightarrow+\infty} \Psi(t)=+\infty$. Suppose that $\phi: \mathbb{R}^{d} \longrightarrow \mathcal{E}$ is a continuous function with values on the separable Banach space $(\mathcal{E},|\cdot|)$. Denote by $\mathrm{B}_{r}$ the open ball in $\mathbb{R}^{d}$ centered at 0 with radius $r$. Then, provided

$$
\begin{equation*}
\Gamma=\int_{\mathrm{B}_{r}} \int_{\mathrm{B}_{r}} \Psi\left(\frac{|\phi(x)-\phi(y)|}{\rho(|x-y|)}\right) \mathrm{d} x \mathrm{~d} y<+\infty \tag{4}
\end{equation*}
$$

it holds, for all $x, y \in \mathrm{~B}_{r}$,

$$
\begin{equation*}
|\phi(x)-\phi(y)| \leq 8 \int_{0}^{2|x-y|} \Psi^{-1}\left(\frac{4^{d+1} \Gamma}{\lambda_{d} u^{2 d}}\right) \rho(\mathrm{d} u) \tag{5}
\end{equation*}
$$

where $\lambda_{d}$ is a universal constant depending only on $d$.
We now aim at proving the following result, which allows to go from pointwise estimates to locally uniform estimates, by assuming Hölder regularity in $\mathbf{L}_{p}$. It will help to check (H2) using much easier conditions.

Theorem 2. Let $p>d$. Assume that $G$ is $\mathcal{F} \otimes \mathcal{B}\left(\mathbb{R}^{d}\right)$-measurable mapping $(\omega, x) \in$ $\left(\Omega, \mathbb{R}^{d}\right) \mapsto G(\omega, x) \in \mathcal{E}$, continuous in $x$ for a.e. $\omega$. Assume that $G(x)$ is in $\mathbf{L}_{p}$ for any $x$ and that there exist constants $C^{(G)} \in[0,+\infty), \beta^{(G)} \in(d / p, 1]$ and $\tau^{(G)} \in$ $[0,+\infty)$ such that

$$
\begin{equation*}
\|G(x)-G(y)\|_{\mathbf{L}_{p}} \leq C^{(G)}|x-y|^{\beta^{(G)}}(1+|x|+|y|)^{\tau^{(G)}}, \quad \forall(x, y) \in \mathbb{R}^{d} \times \mathbb{R}^{d} \tag{6}
\end{equation*}
$$

Then, for any $\beta \in\left(0, \beta^{(G)}-d / p\right)$, we have

$$
\begin{equation*}
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{|G(y)-G(x)|}{|y-x|^{\beta}}\right\|_{\mathbf{L}_{p}} \leq c_{(7)} C^{(G)} \lambda^{\tau^{(G)}+\beta^{(G)}-\beta}, \quad \forall \lambda \geq 1 \tag{7}
\end{equation*}
$$

where $c_{(7)}$ is a constant depending only on $d, p, \beta, \beta^{(G)}, \tau^{(G)}$.
A similar result is proved in [RY99, Theorem 2.1, p.26] using the Kolmogorov criterion, with $x$ and $y$ in a compact set, i.e. with $\tau^{(G)}=0$; the quoted result is not sufficient for our study.

Proof. Since $x \mapsto G(x)$ is a.s. continuous, we can apply Lemma 1 by taking $\Psi(t):=$ $t^{p}$ and $\rho(u):=u^{\gamma_{2}}$ with $\gamma_{2}:=\beta+2 d / p$ : Defining $\Gamma$ as in (4) with $G$ instead of $\phi$, we obtain

$$
\begin{align*}
\mathbb{E}(\Gamma) & =\int_{\mathrm{B}_{r}} \int_{\mathrm{B}_{r}} \mathbb{E}\left(\frac{|G(x)-G(y)|^{p}}{|x-y|^{p \gamma_{2}}}\right) \mathrm{d} x \mathrm{~d} y \\
& \leq\left[C^{(G)}\right]^{p}(1+2|r|)^{p \tau^{(G)}} \int_{\mathrm{B}_{r}} \int_{\mathrm{B}_{r}}|x-y|^{p \beta^{(G)}-p \gamma_{2}} \mathrm{~d} x \mathrm{~d} y \\
& =\left[C^{(G)}\right]^{p}(1+2|r|)^{p \tau^{(G)}} r^{p \beta^{(G)}-p \gamma_{2}+2 d} V_{1} \\
& =\left[C^{(G)}\right]^{p}(1+2|r|)^{p \tau^{(G)}} r^{p\left(\beta^{(G)}-\beta\right)} V_{1} \tag{8}
\end{align*}
$$

where $V_{1}:=\int_{\mathrm{B}_{1}} \int_{\mathrm{B}_{1}}|x-y|^{p \beta^{(G)}-p \gamma_{2}} \mathrm{~d} x \mathrm{~d} y$ is a finite integral since $p \beta^{(G)}-p \gamma_{2}=$ $p\left(\beta^{(G)}-\beta\right)-2 d>-d \Leftrightarrow \beta<\beta^{(G)}-d / p$. This proves that $\mathbb{E}(\Gamma)<+\infty$ thus $\Gamma$ is finite a.s.

Moreover, a direct computation shows that

$$
\int_{0}^{r} \Psi^{-1}\left(\frac{4^{d+1} \Gamma}{\lambda_{d} u^{2 d}}\right) \rho(\mathrm{d} u)=\int_{0}^{r}\left(\frac{4^{d+1} \Gamma}{\lambda_{d} u^{2 d}}\right)^{1 / p} \gamma_{2} u^{\gamma_{2}-1} \mathrm{~d} u=\left(\frac{4^{d+1} \Gamma}{\lambda_{d}}\right)^{1 / p} \frac{\beta+2 d / p}{\beta} r^{\beta}, \quad r \geq 0
$$

Therefore, from the above and (5) we derive

$$
\begin{equation*}
|G(x)-G(y)| \leq 8\left(\frac{4^{d+1}}{\lambda_{d}}\right)^{1 / p} \frac{2^{\beta}(\beta+2 d / p)}{\beta} \Gamma^{1 / p}|y-x|^{\beta} \tag{9}
\end{equation*}
$$

for any $x, y$ with $|x| \leq r$ and $|y| \leq r$. Owing to (8) this implies

$$
\mathbb{E}\left(\left|\sup _{x \neq y,|x| \leq r,|y| \leq r} \frac{|G(y)-G(x)|}{|y-x|^{\beta}}\right|^{p}\right)
$$

$$
\begin{aligned}
& \leq\left[8\left(\frac{4^{d+1}}{\lambda_{d}}\right)^{1 / p} \frac{2^{\beta}(\beta+2 d / p)}{\beta}\right]^{p} \mathbb{E}(\Gamma) \\
& \leq\left[8\left(\frac{4^{d+1}}{\lambda_{d}}\right)^{1 / p} \frac{2^{\beta}(\beta+2 d / p)}{\beta} C^{(G)}(1+2|r|)^{\tau^{(G)}} r^{\left(\beta^{(G)}-\beta\right)}\right]^{p} V_{1} .
\end{aligned}
$$

The proof is complete. Observe that in the proof (see inequality (9)) we more precisely show the a.s. Hölder estimate on $\sup _{x \neq y,|x| \leq r,|y| \leq r} \frac{|G(y)-G(x)|}{|y-x|^{\beta}}$ : This is interesting on its own.

As a consequence, we obtain the following result that may serve to easily check (H1).

Corollary 1. Let consider the assumptions and notations of Theorem 2. Then we have

$$
\begin{equation*}
\left\|\sup _{|x| \leq \lambda}|G(x)|\right\|_{\mathbf{L}_{p}} \leq c_{(10)} \lambda^{\tau^{(G)}+\beta^{(G)}}, \quad \forall \lambda \geq 1 \tag{10}
\end{equation*}
$$

where $c_{(10)}:=\|G(0)\|_{\mathbf{L}_{p}}+c_{(7)} C^{(G)}$ where $c_{(7)}$ is defined in Theorem 2 with $\beta=$ $\left(\beta^{(G)}-d / p\right) / 2$. In particular, the constant $c_{(7)}$ depends only on $d, p, \beta^{(G)}, \tau^{(G)}$.

Proof. Using easy inequalities and applying (7) with $y=0$ and $\beta$ as announced, it readily follows

$$
\begin{aligned}
\left\|\sup _{|x| \leq \lambda}|G(x)|\right\|_{\mathbf{L}_{p}} & \leq\|G(0)\|_{\mathbf{L}_{p}}+\left\|\sup _{0<|x| \leq \lambda} \frac{|G(x)-G(0)|}{|x|^{\beta}}\right\|_{\mathbf{L}_{p}} \lambda^{\beta} \\
& \leq\|G(0)\|_{\mathbf{L}_{p}}+c_{(7)} C^{(G)} \lambda^{\tau^{(G)}+\beta^{(G)}}
\end{aligned}
$$

Since $\lambda \geq 1$ and $\tau^{(G)}+\beta^{(G)} \geq 0$, the proof is complete.

## 3 Application to compound Euler schemes

In this section, let $T$ be a positive and finite time horizon and let us consider a standard filtered probability space $(\Omega, \mathbb{F}, \mathbb{P})$ supporting two $q$-dimensional standard Brownian motions $W=\left(W^{1}, \ldots, W^{q}\right)$ and $B=\left(B^{1}, \ldots, B^{q}\right)$ on $[0, T]$. We consider two $\mathbb{R}^{d}$-valued stochastic processes $X$ and $Y$, solutions of the following stochastic differential equations (SDE for short)

$$
\begin{align*}
& \operatorname{SDE}(\mu, \sigma, W): \mathrm{d} X_{t}(x)=\mu\left(t, X_{t}(x)\right) \mathrm{d} t+\sum_{i=1}^{q} \sigma_{i}\left(t, X_{t}(x)\right) \mathrm{d} W_{t}^{i}, \quad X_{0}(x)=x  \tag{11}\\
& \operatorname{SDE}(b, \gamma, B): \mathrm{d} Y_{t}(y)=b\left(t, Y_{t}(y)\right) \mathrm{d} t+\sum_{i=1}^{q} \gamma_{i}\left(t, Y_{t}(y)\right) \mathrm{d} B_{t}^{i}, \quad Y_{0}(y)=y \tag{12}
\end{align*}
$$

where $\mu, b, \sigma_{i}, \gamma_{i}$ are functions from $[0, T] \times \mathbb{R}^{d}$ into $\mathbb{R}^{d}$, globally Lipschitz in space to ensure the existence of a unique strong solution. Depending on the potential applications, we may require that $B$ and $W$ are the same, or different. Denote by $X_{T}^{N}(x)$ (resp. $\left.Y_{T}^{N}(y)\right)$ the Euler scheme with time step $T / N$ of $X_{T}(x)$ (resp. $\left.Y_{T}(y)\right)$ : Using previous results, we aim at establishing a new convergence result of
the compound scheme $X_{t}^{N}\left(Y_{t}^{N}(y)\right)$ to the compound $\operatorname{SDE} X_{t}\left(Y_{t}(y)\right)$ as $N$ goes to infinity, under the form

$$
\left\|X_{t}^{N}\left(Y_{t}^{N}(y)\right)-X_{t}\left(Y_{t}(y)\right)\right\|_{\mathbf{L}_{p}}=O\left(N^{-1 / 2}\right)
$$

for any $p>0$. For a rigorous statement under precise assumptions, see Theorem 3. This approximation issue, interesting on its own, is actually motivated by other potential applications we briefly expose and that will be subject of future and deeper investigations.

Relation with approximation of stochastic partial differential equations (SPDEs). This work constitutes a first step in a subject that until now has not been addressed to our knowledge, that is to approach solutions of SPDEs by approximating compound SDEs. Relating compound SDEs to SPDEs is, in a sense, obvious since it is sufficient to apply the Itô-Ventzel formula [Kun97, Section 3.3] (under good regularity assumptions on $(\mu, \sigma)$ ) to the compound process $U(t, y):=X_{t}\left(Y_{t}(y)\right)$ to show that $(t, y) \mapsto X_{t}\left(Y_{t}(y)\right)$ solves the second order SPDE, with stochastic coefficients, given by (to simplify we take $d=q=1$ and $W=B$ )

$$
\begin{aligned}
\mathrm{d} U(t, y)= & \left(\partial_{y} U(t, y) \frac{b\left(t, Y_{t}(y)\right)}{\partial_{y} Y(t, y)}+\frac{1}{2}\left(\partial_{y}^{2} U(t, y)-\partial_{y} U(t, y) \frac{\partial_{y}^{2} Y_{t}(y)}{\partial_{y} Y_{t}(y)}\right) \frac{\gamma^{2}\left(t, Y_{t}(y)\right.}{\left(\partial_{y} Y_{t}(y)\right)^{2}}\right. \\
& \left.+\mu(t, U(t, y))+\partial_{x} \sigma(t, U(t, y)) \gamma\left(t, Y_{t}(y)\right)\right) \mathrm{d} t \\
& +\left(\partial_{y} U(t, y) \frac{\gamma\left(t, Y_{t}(y)\right)}{\partial_{y} Y(t, y)}+\sigma(t, U(t, y))\right) \mathrm{d} W_{t}
\end{aligned}
$$

In the reverse direction, i.e. starting from a SPDE, it is more delicate to establish a link with SDEs. But in the recent work [EM13] based on the theory of stochastic flows, El Karoui and Mrad have established a direct connection between a certain utility SPDE and two SDEs. Indeed, being concerned with progressive stochastic utilities $\left(U(t, x): t \geq 0, x \in \mathbb{R}^{d}\right)$ (a.k.a. Forward Utilities or performance processes, see [MZ10]), the authors show that $U$ (under some regularity assumptions) are inevitably solution of a second order fully nonlinear SPDE. Moreover the marginal utility $\partial_{x} U$ is characterized by two SDEs $X$ and $Y$ under the form $U_{x}=X\left(Y^{-1}\right)$. Here $Y^{-1}$ is the inverse flow of $Y$ and can be interpreted as another SDE, see the above reference for details. The current work paves the way to the derivation of convergent approximation of SPDEs of this form.

### 3.1 Hypotheses

We first study approximations on $X$ and for this, we state related assumptions on the $\mathbb{R}^{d}$-valued drift coefficient $\mu=\left\{\mu(t, x) ; t \in[0, T], x \in \mathbb{R}^{d}\right\}$ and the $\mathbb{R}^{d} \times \mathbb{R}^{q}$-valued diffusion coefficient $\sigma=\left\{\sigma_{i}(t, x) ; 1 \leq i \leq q, t \in[0, T], x \in \mathbb{R}^{d}\right\}$ which we suppose to be regular enough in time and space. When we will discuss on approximation of $X(Y)$, similar assumptions will be made on the coefficients $b$ and $\gamma_{i}$ of Equation (12) for $Y$.
(HP1) The coefficients $\mu$ and $\sigma$ are Lipschitz continuous in space uniformly in time. More precisely, there exists a finite constant $C^{X}$ such that for any $t \in[0, T]$ and $x, y \in \mathbb{R}^{d}$

$$
\begin{cases}|\mu(t, x)-\mu(t, y)| \leq C^{X}|x-y|, & |\mu(t, 0)| \leq C^{X}  \tag{HP1}\\ |\sigma(t, x)-\sigma(t, y)| \leq C^{X}|x-y|, & |\sigma(t, 0]| \leq C^{X}\end{cases}
$$

(HP2) $\mu$ and $\sigma$ are continuously space-differentiable functions such that their derivatives $\nabla_{x} \mu:=\left\{\nabla_{x} \mu(t, x) ; t \in[0, T], x \in \mathbb{R}^{d}\right\}$ and $\nabla_{x} \sigma=\left\{\nabla_{x} \sigma_{i}(t, x) ;\right.$ $\left.1 \leq i \leq q, t \in[0, T], x \in \mathbb{R}^{d}\right\}$ are $\delta$-Hölder for a certain exponent $\delta \in(0,1]$. Namely, there exists a finite constant $C^{X, \nabla}$ such that for any $t \in[0, T]$ and $x, y \in \mathbb{R}^{d}$

$$
\begin{cases}\left|\nabla_{x} \mu(t, x)-\nabla_{x} \mu(t, y)\right| \leq C^{X, \nabla}|x-y|^{\delta}, & \left|\nabla_{x} \mu(t, x)\right| \leq C^{X, \nabla} \\ \left|\nabla_{x} \sigma(t, x)-\nabla_{x} \sigma(t, y)\right| \leq C^{X, \nabla}|x-y|^{\delta}, & \left|\nabla_{x} \sigma(t, x)\right| \leq C^{X, \nabla}\end{cases}
$$

(HP2)
(HP3) $\mu$ and $\sigma$ are Hölder continuous in time, locally in space, i.e. there exists an exponent $\alpha \in(0,1]$ and a finite constant $C^{X}$, such that for any $x \in \mathbb{R}^{d}$ and $s, t \in[0, T]$

$$
\begin{equation*}
|\mu(t, x)-\mu(s, x)|+|\sigma(t, x)-\sigma(s, x)| \leq C^{X}(1+|x|)|t-s|^{\alpha} . \tag{HP3}
\end{equation*}
$$

(HP4) $\mu$ and $\sigma$ are continuously space-differentiable functions such that their derivatives are Hölder continuous in time, locally in space, i.e. there exists an exponent $\alpha \in(0,1]$ and a finite constant $C^{X, \nabla}$, such that for any $x \in \mathbb{R}^{d}$ and $s, t \in[0, T]$

$$
\begin{equation*}
\left|\nabla_{x} \mu(t, x)-\nabla_{x} \mu(s, x)\right|+\left|\nabla_{x} \sigma(t, x)-\nabla_{x} \sigma(s, x)\right| \leq C^{X, \nabla}(1+|x|)|t-s|^{\alpha} . \tag{HP4}
\end{equation*}
$$

Denoting in the same way the constants of (HP1) and (HP3) (resp. (HP2) and (HP4)) by $C^{X}$ (resp. $C^{X, \nabla}$ ) is made for the sake of simplicity.
Assumption (HP1) ensures the existence of a strong continuous solution to the $\operatorname{SDE}(\mu, \sigma, W)$, which is adapted to the natural filtration of $W$ completed by the $\mathbb{P}$-null sets: (HP1) plays a crucial role to establish a $\mathbf{L}_{p}$-estimates. It is also wellknown [Kun97, Theorem 4.5.1] that the map $(t, x) \mapsto X_{t}(\omega, x)$ has a modification which is continuous a.s., we shall systematically work with this modification from now on. Assumption (HP2) is a sufficient condition (see [Kun97, Theorem 3.3.3]) under which the above map is $C^{1}$ in $x$. Assumptions (HP3) and (HP4) enable us, essentially, to establish convergence results of the Euler discretization scheme within the paper setting.

### 3.2 Compound Euler schemes: Main result

Under (HP1), let us consider the strong solution to (11): its Euler scheme with $N \geq 1$ discretization times and step-size $\frac{T}{N}$ is defined as usually as follows.

- $\operatorname{Set} X_{0}^{N}(x)=x$.
- For $k=0, \ldots, N-1$ and $t \in\left(k \frac{T}{N},(k+1) \frac{T}{N}\right]$, set

$$
X_{t}^{N}(x)=X_{k \frac{T}{N}}^{N}(x)+\mu\left(k \frac{T}{N}, X_{k \frac{T}{N}}^{N}(x)\right)\left(t-\frac{k T}{n}\right)+\sum_{i=1}^{q} \sigma_{i}\left(k \frac{T}{N}, X_{k \frac{T}{N}}^{N}(x)\right)\left(W_{t}^{i}-W_{k \frac{T}{N}}^{i}\right) .
$$

It can be equivalently written as a continuous Itô process: Denoting by $\tau_{t}:=$ $\left[\frac{N t}{T}\right] \frac{T}{N}$ the last discretization-time before $t$, we have

$$
\begin{equation*}
X_{t}^{N}(x)=x+\int_{0}^{t} \mu\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right) \mathrm{d} s+\sum_{i=1}^{q} \int_{0}^{t} \sigma_{i}\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right) \mathrm{d} W_{s}^{i} . \tag{13}
\end{equation*}
$$

Similarly, assume that $b$ and $\gamma$ fulfills (HP1), so that the strong solution $Y$ to (12) is well defined, together with its Euler scheme $Y^{N}$.

The section is devoted to establish the following main result.
Theorem 3. Assume that $\mu$ and $\sigma$ satisfy Assumptions (HP1), (HP2), (HP3) and (HP4) (which $\alpha$-parameter is denoted by $\alpha^{X}$ ) and that b and $\gamma$ satisfy Assumptions (HP1) and (HP3) (which $\alpha$-parameter is denoted by $\alpha^{Y}$ ).

Then the compound Euler scheme $X^{N}\left(Y^{N}\right)$ converges to $X .(Y$.$) in any \mathbf{L}_{p}$-norm, at the order $\beta:=\min \left(\alpha^{X}, \alpha^{Y}, \frac{1}{2}\right)$ w.r.t. $N$ : For any $p>0$, there is a finite constant $C_{p}$ such that for any $s, t \in[0, T]$

$$
\left\|X_{t}^{N}\left(Y_{s}^{N}\right)-X_{t}\left(Y_{s}\right)\right\|_{\mathbf{L}_{p}} \leq C_{p} N^{-\beta}, \quad \forall N \geq 1
$$

The rest of this section is devoted to its long proof, which requires intermediate estimates on the SDE and its Euler scheme, some of them being completely new (Theorems 4 and 7).

### 3.3 Proof of Theorem 3

In this subsection, we will make use of different constants that may depend on the integer $p$ of $\mathbf{L}_{p}$-norm, on the dimensions $d$ and $q$, on the time horizon $T$ and on the constants from the assumptions: These constants will be called generic constant and will be denoted by the same notation $C_{p}$ even if their values change from line to line. They will not depend on $N$.

We denote by $C_{p}^{\mathrm{BDG}}$ the constant of the upper Burkholder-Davis-Gundy inequality with $\mathbf{L}_{p}$-norm (see the right-hand side of [RY99, Theorem 4.1, p.160]).

### 3.3.1 SDE: differentiability, local and uniform estimates

To analyze the approximation of the compound SDE $X(Y)$, precise estimates on the maps $x \mapsto X_{t}(\omega, x)$ are needed: Such random fields are also called stochastic flows and are the main subject of Kunita's book [Kun97]. As aforementioned, under (HP1), the map $(t, x) \mapsto X_{t}(\omega, x)$ has a continuous modification we are working with. The additional space regularity is connected to the regularity of the coefficients $(\mu, \sigma)$, owing to (HP2), which can be described as follows.

Proposition 2 ([Kun97, Theorem 4.6.5]). Under Assumptions (HP1) and (HP2), the strong solution $X_{t}(x)$ to (11) is continuously differentiable in space and its derivative denoted by $\nabla X_{t}(x)$ is locally $\varepsilon$-Hölder ${ }^{1}$ for any $\varepsilon<\delta$. Furthermore, it is a semimartingale solution of a linear equation, with bounded stochastic parameters $\left(\nabla_{x} \mu\left(t, X_{t}(x)\right), \nabla_{x} \sigma\left(t, X_{t}(x)\right)\right)$ given by

$$
\begin{align*}
\nabla X_{0}(x) & =\mathrm{Id}, \\
\mathrm{~d} \nabla X_{t}(x) & =\nabla_{x} \mu\left(t, X_{t}(x)\right) \nabla X_{t}(x) \mathrm{d} t+\sum_{i=1}^{q} \nabla_{x} \sigma_{i}\left(t, X_{t}(x)\right) \nabla X_{t}(x) \mathrm{d} W_{t}^{i} . \tag{14}
\end{align*}
$$

We now proceed to $L_{p}$-estimates of $X_{t}(x)$ and its sensitivity w.r.t. $x$. We collect several useful results in the following Proposition.

Proposition 3. Assume (HP1). For any $p>0$, there exist generic constants $C_{p,(15)}$ and $C_{p,(16)}$ such that

$$
\begin{align*}
\left\|X_{t}(x)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(15)}(1+|x|),  \tag{15}\\
\left\|X_{t}(x)-X_{t}(y)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(16)}|x-y| \tag{16}
\end{align*}
$$

for any $(t, x, y) \in[0, T] \times \mathbb{R}^{d} \times \mathbb{R}^{d}$. In addition under (HP2), for any $p>0$ there exist generic constants $C_{p,(17)}$ and $C_{p,(18)}$ such that

$$
\begin{align*}
\left\|\nabla X_{t}(x)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(17)},  \tag{17}\\
\left\|\nabla X_{t}(x)-\nabla X_{t}(y)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(18)}|x-y|^{\delta} \tag{18}
\end{align*}
$$

for any $(t, x, y) \in[0, T] \times \mathbb{R}^{d} \times \mathbb{R}^{d}$.
Proof. The proofs of inequalities (15) and (16) are standard, see [Kun97, Lemmas 4.5.3 and 4.5.5]. The uniform estimate (17) is also easy to obtain, in view of (14) and owing to the boundedness of $\nabla_{x} \mu$ and $\nabla_{x} \sigma_{i}$, we leave the details to the reader.

It remains to show (18) under (HP2). To alleviate the notation, we provide the proof when $d=q=1$, the general case being similar. Also, we can focus on the case $p \geq 2$ since we can deduce the result for $p<2$ using the stability of $\mathbf{L}_{p}$-norm combined with the result for $p=2$. First, from (14) write

$$
\begin{aligned}
\nabla X_{t}(x)-\nabla X_{t}(y) & =\int_{0}^{t} \nabla_{x} \mu\left(s, X_{s}(x)\right)\left(\nabla X_{s}(x)-\nabla X_{s}(y)\right) \mathrm{d} s \\
& +\int_{0}^{t}\left(\nabla_{x} \mu\left(s, X_{s}(s)\right)-\nabla_{x} \mu\left(s, X_{s}(y)\right)\right) \nabla X_{s}(y) \mathrm{d} s \\
& +\int_{0}^{t} \nabla_{x} \sigma\left(s, X_{s}(x)\right)\left(\nabla X_{s}(x)-\nabla X_{s}(y)\right) \mathrm{d} W_{s} \\
& +\int_{0}^{t}\left(\nabla_{x} \sigma\left(s, X_{s}(s)\right)-\nabla_{x} \sigma\left(s, X_{s}(y)\right)\right) \nabla X_{s}(y) \mathrm{d} W_{s} .
\end{aligned}
$$

[^1]Take the power $p$ and the expectation, then apply the Burkholder-Davis-Gundy inequality, the Jensen equality ( $p \geq 2$ ) and the Cauchy-Schwarz inequality; it leads to
$\mathbb{E}\left(\left|\nabla X_{t}(x)-\nabla X_{t}(y)\right|^{p}\right)$
$\leq 4^{p-1} t^{p-1} \int_{0}^{t} \mathbb{E}\left(\left|\nabla_{x} \mu\left(s, X_{s}(x)\right)\left(\nabla X_{s}(x)-\nabla X_{s}(y)\right)\right|^{p}\right) \mathrm{d} s$
$+4^{p-1} t^{p-1} \int_{0}^{t} \sqrt{\mathbb{E}\left(\left|\nabla_{x} \mu\left(s, X_{s}(s)\right)-\nabla_{x} \mu\left(s, X_{s}(y)\right)\right|^{2 p}\right)} \sqrt{\mathbb{E}\left(\left|\nabla X_{s}(y)\right|^{2 p}\right)} \mathrm{d} s$
$+4^{p-1}\left[C_{p}^{\mathrm{BDG}}\right]^{p} t^{p / 2-1} \int_{0}^{t} \mathbb{E}\left(\left|\nabla_{x} \sigma\left(s, X_{s}(x)\right)\left(\nabla X_{s}(x)-\nabla X_{s}(y)\right)\right|^{p}\right) \mathrm{d} s$
$+4^{p-1}\left[C_{p}^{\mathrm{BDG}}\right]^{p} t^{p / 2-1} \int_{0}^{t} \sqrt{\mathbb{E}\left(\left|\nabla_{x} \sigma\left(s, X_{s}(s)\right)-\nabla_{x} \sigma\left(s, X_{s}(y)\right)\right|^{2 p}\right)} \sqrt{\mathbb{E}\left(\left|\nabla X_{s}(y)\right|^{2 p}\right)} \mathrm{d} s$.
Now, take advantage of the Assumptions (HP1) and (HP2), together with the estimates (16) and (17): it readily follows that $\iota(t):=\mathbb{E}\left(\left|\nabla X_{t}(x)-\nabla X_{t}(y)\right|^{p}\right)$ solves

$$
\begin{aligned}
\iota(t) \leq & 4^{p-1}\left[C^{X, \nabla}\right]^{p}\left(T^{p-1}+\left[C^{\mathrm{BDG}}\right]^{p} T^{p / 2-1}\right) \int_{0}^{t} \iota(s) \mathrm{d} s \\
& +4^{p-1}\left[C^{X, \nabla}\right]^{p} C_{2 p,(17)}^{p} C_{2 p \delta,(16)}^{p \delta}\left(T^{p}+\left[C^{\mathrm{BDG}}\right]^{p} T^{p / 2}\right)|x-y|^{p \delta}
\end{aligned}
$$

The estimate (18) is then a direct consequence of Gronwall's lemma.

Thanks to the results of Section 2, we are now in a position to generalize Proposition 3 by putting the sup over the space variable inside the expectation. This is the following assertion, which is a new result to our knowledge.

Theorem 4. Assume Assumption (HP1). For any $p>0$ and any $\beta \in(0,1)$, there exist generic constants $C_{p,(19)}$ and $C_{p,(20)}$ such that, for any $t \in[0, T]$,

$$
\begin{gather*}
\left\|\sup _{|x| \leq \lambda}\left|X_{t}(x)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(19)} \lambda, \quad \forall \lambda \geq 1,  \tag{19}\\
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{\left|X_{t}(x)-X_{t}(y)\right|}{|y-x|^{\beta}}\right\|_{\mathbf{L}_{p}} \leq C_{p,(20)} \lambda^{1-\beta}, \quad \forall \lambda \geq 1 . \tag{20}
\end{gather*}
$$

Assume furthermore Assumption (HP2). For any $p>0$ and any $\beta \in(0, \delta)$, there exist generic constants $C_{p,(21)}, C_{p,(22)}$ and $C_{p,(23)}$ such that, for any $t \in[0, T]$,

$$
\begin{array}{r}
\left\|\sup _{|x| \leq \lambda}\left|\nabla X_{t}(x)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(21)} \lambda^{\delta}, \quad \forall \lambda \geq 1, \\
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{\left|\nabla X_{t}(x)-\nabla X_{t}(y)\right|}{|y-x|^{\beta}}\right\|_{\mathbf{L}_{p}} \leq C_{p,(22)} \lambda^{\delta-\beta}, \quad \forall \lambda \geq 1, \\
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{\left|X_{t}(x)-X_{t}(y)\right|}{|y-x|}\right\|_{\mathbf{L}_{p}} \leq C_{p,(23)} \lambda^{\delta}, \quad \forall \lambda \geq 1 \tag{23}
\end{array}
$$

Proof. Let $\beta \in(0,1)$ : We first show (19) and (20) for any $p>d /(1-\beta)>d$. Owing to (16), we can apply Theorem 2 to $G(x):=X_{t}(x)$ with $\beta^{(G)}=1 \in(d / p, 1]$ and $\tau^{(G)}=0$, to conclude that (20) holds with the given index $\beta$ since $\beta<1-d / p \Leftrightarrow p>$ $d /(1-\beta)$. Moreover the application of Corollary 1 provides (19). Now it remains to relax the constraint on $p$ : for $p \leq d /(1-\beta)$, set $\bar{p}=2 d /(1-\beta)$ for which (20) holds and take advantage of the stability property of $\mathbf{L}_{p}$-norm to write

$$
\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{\left|X_{t}(x)-X_{t}(y)\right|}{|y-x|^{\beta}}\right\|_{\mathbf{L}_{p}} \leq\left\|\sup _{x \neq y,|x| \leq \lambda,|y| \leq \lambda} \frac{\left|X_{t}(x)-X_{t}(y)\right|}{|y-x|^{\beta}}\right\|_{\mathbf{L}_{\bar{p}}} \leq C_{\bar{p},(20)} \lambda^{1-\beta}
$$

The same arguments apply to prove that (19) holds for any $p>0$.
The justification of (21) and (22) follows the same arguments as above, using (18) instead of (16): then Theorem 2 and Corollary 1 can be applied to $G(x):=\nabla X_{t}(x)$ with $\beta^{(G)}=\delta$ and $\tau^{(G)}=0$. We leave the details to the reader.

Last, observe that for any $x, y$ such that $|x| \leq \lambda$ and $|y| \leq \lambda$, we have $\mid X_{t}(x)-$ $X_{t}(y)\left|\leq \sup _{|z| \leq \lambda}\right| \nabla X_{t}(z)| | y-x \mid:$ thus, (23) readily follows from (21).

Observe that the additional smoothness in (HP2) enables us to improve (20) (for $\beta<1$ ) to (23) (i.e. $\beta=1$ ): this improvement will play an important role in the derivation of Theorem 3.

### 3.3.2 Euler scheme: local and uniform estimates

Still as intermediate steps to prove Theorem 3, we partly generalize the previous results about the SDE to its Euler approximation. Some derivations are more subtle and require details at some places. Recall the definition of Euler scheme in (13).

First, as for the solution of the $\operatorname{SDE}(\mu, \sigma)$, some estimates for its approximation scheme are needed. This is the analogue of Proposition 3.

Proposition 4. Under (HP1), for any $p>0$ there exist generic constants $C_{p,(24)}$ and $C_{p,(25)}$ such that

$$
\begin{align*}
\left\|X_{t}^{N}(x)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(24)}(1+|x|)  \tag{24}\\
\left\|X_{t}^{N}(x)-X_{t}^{N}(y)\right\|_{\mathbf{L}_{p}} & \leq C_{p,(25)}|x-y| \tag{25}
\end{align*}
$$

for any $(t, x, y) \in[0, T] \times \mathbb{R}^{d} \times \mathbb{R}^{d}$.
We omit the proof which is quite standard. Following the same arguments than for the SDE case (Theorem 4), we can put the sup over the space variable inside the $\mathbf{L}_{p}$-norm, it gives the following.

Proposition 5. Under (HP1), the estimates (19) and (20) where we replace $X$ by $X^{N}$ hold true, up to changing the generic constants.

Let us now show the following estimates on local increments, it will be needed for the sequel.

Lemma 2. Assume Assumption (HP1) and let $p>0$. Then there exist generic constants $C_{p,(26)}$ and $C_{p,(27)}$ such that, for any $x, y \in \mathbb{R}^{d}$ and any $t \in[0, T]$,

$$
\begin{align*}
&\left\|\sup _{\tau_{t} \leq u \leq t}\left|X_{u}^{N}(x)-X_{\tau_{u}}^{N}(x)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(26)} \frac{(1+|x|)}{N^{1 / 2}},  \tag{26}\\
&\left\|\sup _{\tau_{t} \leq u \leq t}\left|X_{u}^{N}(x)-X_{u}^{N}(y)-X_{\tau_{u}}^{N}(x)+X_{\tau_{u}}^{N}(y)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(27)} \frac{|x-y|}{N^{1 / 2}} . \tag{27}
\end{align*}
$$

Proof. Here again, it is enough to prove the estimates for $p \geq 2$, which we assume from now on. Also we take $d=q=1$ to simplify the exposure. Similarly to the proof of Proposition 3, Burkholder-Davis-Gundy's inequality combined with Jensen's inequality readily leads to

$$
\begin{aligned}
\mathbb{E}\left(\sup _{\tau_{t} \leq u \leq t}\left|X_{u}^{N}(x)-X_{\tau_{u}}^{N}(x)\right|^{p}\right) & \leq 2^{p-1}\left(\left(t-\tau_{t}\right)^{p-1} \int_{\tau_{t}}^{t} \mathbb{E}\left(\left|\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)\right|^{p}\right) \mathrm{d} s\right. \\
& \left.+\left(t-\tau_{t}\right)^{p / 2-1}\left[C_{p}^{\mathrm{BDG}}\right]^{p} \int_{\tau_{t}}^{t} \mathbb{E}\left(\left|\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)\right|^{p}\right) \mathrm{d} s\right) .
\end{aligned}
$$

Finally, from Assumption (HP1), we have $|\mu(t, x)|+|\sigma(t, x)| \leq C^{X}(1+|x|)$ for any $t \in[0, T]$; combined with (24), we deduce

$$
\begin{aligned}
\mathbb{E}\left(\sup _{\tau_{t} \leq u \leq t}\left|X_{u}^{N}(x)-X_{\tau_{u}}^{N}(x)\right|^{p}\right) & \leq 2^{p-1}\left(C^{X}\right)^{p}\left(\left(t-\tau_{t}\right)^{p} 2^{p-1}\left(1+C_{p,(24)}^{p}(1+|x|)^{p}\right)\right. \\
& \left.+\left[C_{p}^{\mathrm{BDG}}\right]^{p}\left(t-\tau_{t}\right)^{p / 2} 2^{p-1}\left(1+C_{p,(24)}^{p}(1+|x|)^{p}\right)\right)
\end{aligned}
$$

which readily leads to the announced estimate (26).
Let us now turn to the second inequality: The same arguments combined with Assumption (HP1) and (25) lead to

$$
\begin{aligned}
& \mathbb{E}\left(\sup _{\tau_{t} \leq u \leq t}\left|X_{u}^{N}(x)-X_{u}^{N}(y)-X_{\tau_{u}}^{N}(x)+X_{\tau_{u}}^{N}(y)\right|^{p}\right) \\
& \leq 2^{p-1}\left(t-\tau_{t}\right)^{p-1} \int_{\tau_{t}}^{t} \mathbb{E}\left(\left|\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p}\right) \mathrm{d} s \\
& +2^{p-1}\left[C_{p}^{\mathrm{BDG}}\right]^{p}\left(t-\tau_{t}\right)^{p / 2-1} \int_{\tau_{t}}^{t} \mathbb{E}\left(\left|\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p}\right) \mathrm{d} s \\
& \leq 2^{p-1}\left(C^{X}\right)^{p} C_{p,(25)}^{p}\left(\left(t-\tau_{t}\right)^{p}+\left[C_{p}^{\mathrm{BDG}}\right]^{p}\left(t-\tau_{t}\right)^{p / 2}\right)|x-y|^{p},
\end{aligned}
$$

which completes the proof.

Strong convergence (classical result). Since in the Euler scheme dynamics the coefficients $\mu$ and $\sigma$ are computed at the left of each time interval, we need to account for their time regularity in order to derive a sharp convergence result: This is stated through Assumption (HP3). The proof of the following result can be found in [BL93, Theorem B.1.4 p. 276].

Theorem 5. Assume Assumptions (HP1) and (HP3) and set $\beta=\min \left(\alpha, \frac{1}{2}\right)$. Then, for any $p>0$ there exists a generic constant $C_{p,(28)}$ such that for any $x \in \mathbb{R}^{d}$

$$
\begin{equation*}
\left\|\sup _{t \leq T}\left|X_{t}(x)-X_{t}^{N}(x)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(28)} \frac{(1+|x|)}{N^{\beta}} . \tag{28}
\end{equation*}
$$

Furthermore, for any $\gamma<\beta$, the random variables $\left(N^{\gamma} \sup _{t \leq T}\left|X_{t}-X_{t}^{N}\right|\right)_{N \geq 1}$ converge almost surely to 0 as $N$ tends to $+\infty$.

Unfortunately, the classical estimate of Theorem 5 is not sufficient to analyze the error of compound Euler schemes: in view of Theorem 1 and its assumptions (in particular (H3)), one should have a sup over $|x| \leq \lambda$ inside the $\mathbf{L}_{p}$-norm. This is the purpose of the next derivations.

Strong convergence (new results). To obtain locally uniform in space convergence results, the supplementary assumptions of regularity in space and time for $\nabla_{x} \mu$ and $\nabla_{x} \sigma_{i}($ see (HP2) and (HP4)) are seemingly important. Thus Theorem 5 can be generalized to the following crucial one.

Theorem 6. Assume (HP1), (HP2), (HP3), (HP4) and let $\beta=\min \left(\alpha, \frac{1}{2}\right)$. For any $p>0$, there exists a generic constant $C_{p,(29)}$ such that

$$
\begin{equation*}
\left\|\sup _{u \leq t}\left|X_{u}(x)-X_{u}^{N}(x)-X_{u}(y)+X_{u}^{N}(y)\right|\right\|_{\mathbf{L}_{p}} \leq C_{p,(29)}(1+|x|+|y|) \frac{|x-y|+|x-y|^{\delta}}{N^{\beta}} \tag{29}
\end{equation*}
$$

for all $x, y \in \mathbb{R}^{d}$ and $t \in[0, T]$.
Similarly to Theorem 4, we can now derive estimates locally uniformly in space.
Theorem 7. Under Assumptions of Theorem 6, for any $p>0$ there exists a finite generic constant $C_{p,(30)}$ such that, for any $t \in[0, T]$,

$$
\begin{equation*}
\left\|\sup _{|x| \leq \lambda}\left|X_{t}(x)-X_{t}^{N}(x)\right|\right\|_{\mathbf{L}_{p}} \leq \frac{C_{p,(30)}}{N^{\beta}} \lambda^{2}, \quad \forall \lambda \geq 1 \tag{30}
\end{equation*}
$$

Proof. We aim at applying Corollary 1 by checking the assumptions of Theorem 2 applied to $G(x):=X_{t}(x)-X_{t}^{N}(x)$. From (29) we have

$$
\begin{aligned}
\|G(x)-G(y)\|_{\mathbf{L}_{p}} & \leq C_{p,(29)}(1+|x|+|y|) \frac{|x-y|+|x-y|^{\delta}}{N^{\beta}} \\
& \leq 2 C_{p,(29)}(1+|x|+|y|)^{2-\delta} \frac{|x-y|^{\delta}}{N^{\delta}}
\end{aligned}
$$

using $|x-y|+|x-y|^{\delta}=|x-y|^{\delta}\left(1+|x-y|^{1-\delta}\right) \leq 2|x-y|^{\delta}(1+|x|+|y|)^{1-\delta}$. Thus, we can take $C^{(G)}=2 C_{p,(29)} / N^{\delta}, \tau^{(G)}=2-\delta$ and $\beta^{(G)}=\delta$ provided that $\delta \in(d / p, 1]$, which is true for $p$ large enough. Therefore for such $p$, the estimate (10) holds true, which is the announced inequality of Theorem 7 . The estimate for smaller values of $p$ are automatically satisfied invoking once again the stability of $\mathbf{L}_{p}$ norms as $p$ decreases.

Proof of Theorem 6. As in the previous proofs, we argue that it is enough to assume $p \geq 2$. To alleviate the presentation, we additionally assume $d=q=1$, the derivation in the general case being similar. From the dynamics of $X$ and $X^{N}$, we write

$$
X_{t}(x)-X_{t}^{N}(x)-X_{t}(y)+X_{t}^{N}(y)
$$

$$
\begin{aligned}
& =\int_{0}^{t}\left(\mu\left(s, X_{s}(x)\right)-\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\mu\left(s, X_{s}(y)\right)+\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right) \mathrm{d} s \\
& +\int_{0}^{t}\left(\sigma\left(s, X_{s}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right) \mathrm{d} W_{s}
\end{aligned}
$$

Then, as in the proof of Proposition 3, we obtain

$$
\begin{align*}
& \mathbb{E}\left(\sup _{u \leq t}\left|X_{u}(x)-X_{u}^{N}(x)-X_{u}(y)+X_{u}^{N}(y)\right|^{p}\right)  \tag{31}\\
& \leq 2^{p-1} t^{p-1} \int_{0}^{t} \mathbb{E}\left(\left|\mu\left(s, X_{s}(x)\right)-\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\mu\left(s, X_{s}(y)\right)+\mu\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p}\right) \mathrm{d} s \\
& +2^{p-1}\left[C_{p}^{\mathrm{BDG}}\right]^{p} t^{p / 2-1} \int_{0}^{t} \mathbb{E}\left(\left|\sigma\left(s, X_{s}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p}\right) \mathrm{d} s
\end{align*}
$$

Actually, both terms of the right side of above inequality can be treated in the same way, thus we only detail the computations for the second integral. First write that

$$
\begin{align*}
& \sigma\left(s, X_{s}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right) \\
& \quad=\sigma\left(s, X_{s}(x)\right)-\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(s, X_{s}^{N}(y)\right) \\
& +\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}^{N}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right) \tag{32}
\end{align*}
$$

Now, we treat the two lines above separately.
Step 1. Denoting by $X_{s}^{N, \lambda, x}:=X_{s}(x)+\lambda\left(X_{s}^{N}(x)-X_{s}(x)\right)$ for $\lambda \in[0,1]$, we have

$$
\begin{aligned}
& \sigma\left(s, X_{s}(x)\right)-\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(s, X_{s}^{N}(y)\right) \\
& =\left(X_{s}(x)-X_{s}^{N}(x)-X_{s}(y)+X_{s}^{N}(y)\right) \int_{0}^{1} \nabla_{x} \sigma\left(s, X_{s}^{N, \lambda, x}\right) \mathrm{d} \lambda \\
& +\left(X_{s}(y)-X_{s}^{N}(y)\right) \int_{0}^{1}\left(\nabla_{x} \sigma\left(s, X_{s}^{N, \lambda, x}\right)-\nabla_{x} \sigma\left(s, X_{s}^{N, \lambda, y}\right)\right) \mathrm{d} \lambda
\end{aligned}
$$

Now we use the definition of the process $X^{N, \lambda, x}$, the fact that $\left|\nabla_{x} \sigma(t, x)\right| \leq C^{X, \nabla}$ and $\left|\nabla_{x} \sigma(t, x)-\nabla_{x} \sigma(t, y)\right| \leq C^{X, \nabla}|x-y|^{\delta}$; we then deduce (for a generic constant $C_{p}$ which values may change from line to line)

$$
\begin{aligned}
& \left|\sigma\left(s, X_{s}(x)\right)-\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(s, X_{s}^{N}(y)\right)\right|^{p} \\
& \leq C_{p}\left(\left|X_{s}(x)-X_{s}^{N}(x)-X_{s}(y)+X_{s}^{N}(y)\right|^{p}\right. \\
& \left.\quad+\left|X_{s}(y)-X_{s}^{N}(y)\right|^{p} \int_{0}^{1}\left|(1-\lambda)\left(X_{s}(x)-X_{s}(y)\right)+\lambda\left(X_{s}^{N}(x)-X_{s}^{N}(y)\right)\right|^{\delta p} \mathrm{~d} \lambda\right) \\
& \leq C_{p}\left[\left|X_{s}(x)-X_{s}^{N}(x)-X_{s}(y)+X_{s}^{N}(y)\right|^{p}\right. \\
& \left.\quad \quad \quad+\left|X_{s}(y)-X_{s}^{N}(y)\right|^{p}\left(\left|X_{s}(x)-X_{s}(y)\right|^{\delta p}+\left|X_{s}^{N}(x)-X_{s}^{N}(y)\right|^{\delta p}\right)\right]
\end{aligned}
$$

where we have invoked the Minkowsky inequality to handle the $\mathrm{d} \lambda$-integral and also used (2). From this, integrating over $(s, \omega)$ and applying the Cauchy-Schwarz inequality, we obtain (with a larger constant $C_{p}$ )
$\mathbb{E} \int_{0}^{t}\left|\sigma\left(s, X_{s}(x)\right)-\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(s, X_{s}^{N}(y)\right)\right|^{p} \mathrm{~d} s$

$$
\begin{aligned}
& \leq C_{p}\left[\int_{0}^{t} \mathbb{E}\left(\left|X_{s}(x)-X_{s}^{N}(x)-X_{s}(y)+X_{s}^{N}(y)\right|^{p}\right) \mathrm{d} s\right. \\
& \left.+\int_{0}^{t} \sqrt{\mathbb{E}\left(\left|X_{s}(y)-X_{s}^{N}(y)\right|^{2 p}\right)} \sqrt{\mathbb{E}\left(\left|X_{s}(x)-X_{s}(y)\right|^{2 \delta p}+\left|X_{s}^{N}(x)-X_{s}^{N}(y)\right|^{2 \delta p}\right)} \mathrm{d} s\right]
\end{aligned}
$$

which rewrites, owing to (16)-(25) and (28),

$$
\begin{align*}
& \mathbb{E} \int_{0}^{t}\left|\sigma\left(s, X_{s}(x)\right)-\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(s, X_{s}^{N}(y)\right)\right|^{p} \mathrm{~d} s \\
& \left.\leq C_{p}\left(\int_{0}^{t} \mathbb{E}\left(\mid X_{s}(x)\right)-X_{s}^{N}(x)-X_{s}(y)+\left.X_{s}^{N}(y)\right|^{p}\right) \mathrm{d} s+\frac{(1+|y|)^{p}}{N^{\beta p}}|x-y|^{\delta p}\right) \tag{33}
\end{align*}
$$

for a new generic constant $C_{p}$.
Step 2. Now we are concerned by the second line of Identity (32). Similarly to before, we can write

$$
\begin{aligned}
& \sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}^{N}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right) \\
&= \sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(s, X_{\tau_{s}}^{N}(x)\right)-\left(\sigma\left(s, X_{s}^{N}(y)\right)-\sigma\left(s, X_{\tau_{s}}^{N}(y)\right)\right) \\
&+ \sigma\left(s, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{\tau_{s}}^{N}(y)\right)-\left(\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right) \\
&= \int_{0}^{1} \nabla_{x} \sigma\left(s, X_{s}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(x)-X_{s}^{N}(x)\right)\right) \mathrm{d} \lambda\left(X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)\right) \\
&- \int_{0}^{1} \nabla_{x} \sigma\left(s, X_{s}^{N}(y)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{s}^{N}(y)\right)\right) \mathrm{d} \lambda\left(X_{s}^{N}(y)-X_{\tau_{s}}^{N}(y)\right) \\
&+\int_{0}^{1} \nabla_{x} \sigma\left(s, X_{\tau_{s}}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{\tau_{s}}^{N}(x)\right)\right) \mathrm{d} \lambda\left(X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right) \\
&-\int_{0}^{1} \nabla_{x} \sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{\tau_{s}}^{N}(x)\right)\right) \mathrm{d} \lambda\left(X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right) \\
&= \int_{0}^{1}\left(\nabla_{x} \sigma\left(s, X_{s}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(x)-X_{s}^{N}(x)\right)\right)-\nabla_{x} \sigma\left(s, X_{s}^{N}(y)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{s}^{N}(y)\right)\right)\right) \mathrm{d} \lambda \\
& \quad \times\left(X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)\right) \\
&+\int_{0}^{1} \nabla_{x} \sigma\left(s, X_{s}^{N}(y)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{s}^{N}(y)\right)\right) \mathrm{d} \lambda\left(X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)-X_{s}^{N}(y)+X_{\tau_{s}}^{N}(y)\right) \\
&+ \int_{0}^{1}\left(\nabla_{x} \sigma\left(s, X_{\tau_{s}}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{\tau_{s}}^{N}(x)\right)\right)-\nabla_{x} \sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{\tau_{s}}^{N}(x)\right)\right)\right) \mathrm{d} \lambda \\
& \quad \times\left(X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right) .
\end{aligned}
$$

Now, by taking advantage of the boundedness and regularity assumptions on $\nabla_{x} \sigma$, it readily follows

$$
\begin{aligned}
& \left|\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}^{N}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right| \\
& \leq C^{X, \nabla}\left[\int_{0}^{1}\left|(1-\lambda)\left(X_{s}^{N}(x)-X_{s}^{N}(y)\right)+\lambda\left(X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right)\right|^{\delta} \mathrm{d} \lambda\left|X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)\right|\right. \\
& \quad+\left|X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)-X_{s}^{N}(y)+X_{\tau_{s}}^{N}(y)\right| \\
& \left.\quad+\left|s-\tau_{s}\right|^{\alpha} \int_{0}^{1}\left(1+\left|X_{\tau_{s}}^{N}(x)+\lambda\left(X_{\tau_{s}}^{N}(y)-X_{\tau_{s}}^{N}(x)\right)\right|\right) \mathrm{d} \lambda\left|X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right|\right]
\end{aligned}
$$

By taking the power $p$ and integrating w.r.t. $(s, \omega)$, we get, after standard computations,

$$
\begin{aligned}
& \int_{0}^{t} \mathbb{E}\left(\left|\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}^{N}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p}\right) \mathrm{d} s \\
& \leq C_{p}\left[\int_{0}^{t}\left(\sqrt{\mathbb{E}\left(\left|X_{s}^{N}(x)-X_{s}^{N}(y)\right|^{2 p \delta}\right)}+\sqrt{\left.\left.\mathbb{E}\left(\mid X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right)\right|^{2 p \delta}\right)}\right)\right. \\
& \quad \times \sqrt{\mathbb{E}\left(\left|X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)\right|^{2 p}\right)} \mathrm{d} s \\
& \quad+\int_{0}^{t} \mathbb{E}\left(\left|X_{s}^{N}(x)-X_{\tau_{s}}^{N}(x)-X_{s}^{N}(y)+X_{\tau_{s}}^{N}(y)\right|^{p}\right) \mathrm{d} s \\
& \left.+\frac{1}{N^{\alpha p}} \int_{0}^{t}\left(1+\sqrt{\mathbb{E}\left(\left|X_{\tau_{s}}^{N}(x)\right|^{2 p}\right)}+\sqrt{\mathbb{E}\left(\left|X_{\tau_{s}}^{N}(y)\right|^{2 p}\right)}\right) \sqrt{\mathbb{E}\left(\left|X_{\tau_{s}}^{N}(x)-X_{\tau_{s}}^{N}(y)\right|^{2 p}\right)} \mathrm{d} s\right]
\end{aligned}
$$

for some new generic constant $C_{p}$. Finally, by plugging into the above the results of Proposition 4 and Lemma 2, we obtain (for a new constant $C_{p}$ )

$$
\begin{align*}
& \mathbb{E} \int_{0}^{t}\left|\sigma\left(s, X_{s}^{N}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}^{N}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p} \mathrm{~d} s \\
& \leq C_{p}\left(\frac{|x-y|^{p \delta}}{N^{p / 2}}(1+|x|)^{p}+\frac{|x-y|^{p}}{N^{p / 2}}+\frac{|x-y|^{p}}{N^{\alpha p}}\left(1+|x|^{p}+|y|^{p}\right)\right) \\
& \leq C_{p}(1+|x|+|y|)^{p} \frac{|x-y|^{p}+|x-y|^{\delta p}}{N^{\beta p}} \tag{34}
\end{align*}
$$

We then obtain, by combining (32), (33) and (34),

$$
\begin{aligned}
& \mathbb{E} \int_{0}^{t}\left|\sigma\left(s, X_{s}(x)\right)-\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(x)\right)-\sigma\left(s, X_{s}(y)\right)+\sigma\left(\tau_{s}, X_{\tau_{s}}^{N}(y)\right)\right|^{p} \mathrm{~d} s \\
& \leq C_{p}\left[\int_{0}^{t} \mathbb{E}\left(\left|X_{s}(x)-X_{s}^{N}(x)-X_{s}(y)+X_{s}^{N}(y)\right|^{p}\right) \mathrm{d} s\right. \\
& \left.\quad \quad+(1+|x|+|y|)^{p} \frac{|x-y|^{p}+|x-y|^{\delta p}}{N^{\beta p}}\right]
\end{aligned}
$$

for some new constant $C_{p}$. The same estimates hold for $\mu$ instead of $\sigma$. Hence, plugging the above into (31), we obtain the existence of generic constants $C_{p}$ such that

$$
\begin{aligned}
& \mathbb{E}\left(\sup _{u \leq t}\left|X_{u}(x)-X_{u}^{N}(x)-X_{u}(y)+X_{u}^{N}(y)\right|^{p}\right) \\
& \leq C_{p}\left[\int_{0}^{t} \mathbb{E}\left(\sup _{u \leq s}\left|X_{u}(x)-X_{u}^{N}(x)-X_{u}(y)+X_{u}^{N}(y)\right|^{p}\right) \mathrm{d} s\right. \\
& \left.\quad+(1+|x|+|y|)^{p} \frac{|x-y|^{p}+|x-y|^{\delta p}}{N^{\beta p}}\right] \\
& \quad \leq C_{p}(1+|x|+|y|)^{p} \frac{|x-y|^{p}+|x-y|^{\delta p}}{N^{\beta p}}
\end{aligned}
$$

where the last inequality follows from Gronwall's Lemma; the proof is complete.

### 3.3.3 Completion of the proof of Theorem 3

We now aim at applying Theorem 1 with $F(\omega, x):=X_{t}(\omega, x), F^{N}(\omega, x):=X_{t}^{N}(\omega, x)$, $\Theta:=Y_{s}(\omega, y)$ and $\Theta^{N}:=Y_{s}^{N}(\omega, y)$.

Assumption (H1) is satisfied with $C_{p}^{(\mathbf{H 1})}:=C_{p,(19)}$ and $\alpha_{p}^{(\mathbf{H 1 )})}:=1$ in view of Theorem 4.

Thanks to the inequality (23) of Theorem 4, Assumption (H1) holds true with $\kappa:=1, C_{p}^{(\mathbf{H} 2)}:=C_{p,(23)}$ and $\alpha_{p}^{(\mathbf{H} 2)}:=\delta$.

Moreover (H3) is valid owing to Theorem 7 where we take $C_{p}^{N,(\mathbf{H} 3)}:=\frac{C_{p,(30)}}{N^{X}}$ (with $\beta^{X}:=\min \left(\alpha^{X}, \frac{1}{2}\right)$ ) and $\alpha_{p}^{(\mathbf{H 3})}:=2$.

Last, (H4) is clearly true using

- Propositions 3 and 4 applied to $Y$ instead of $X$, which yields $C_{p}^{(\mathbf{H 4 - a})}:=$ $\max \left(C_{p,(15)}, C_{p,(24)}\right)(1+|y|)$,
- Theorem 5, applied to $Y$ and $Y^{N}$, which gives $C_{p}^{N,(\mathbf{H} 4-\mathrm{b})}:=\frac{C_{p,(28)}}{N^{\beta^{Y}}}(1+|y|)$ with $\beta^{Y}:=\min \left(\alpha^{Y}, \frac{1}{2}\right)$.
We are done.


## 4 Application to stochastic processes at random times

Considering stochastic processes at random times is interesting on its own, and besides it has also many applications: among others, we mention the Dambis-DubinsSchwarz representation of martingale as time-changed Brownian motion [RY99, Chapter V], the Skorokhod Embedding Problem to represent any distribution using Brownian motion stopped at a suitable stopping time [Obł04], and the Brownian motion at Brownian time to derive Feynman-Kac formulas for bi-Laplacian PDEs [Fun79].

We illustrate our general result (Theorem 1) in non trivial applications. To be pedagogical, we start with martingale models at stopping times: here, the usual stochastic calculus tools enable to derive error bounds when only the stopping time $\Theta$ is approximated (Proposition 7). Passing to arbitrary random times is not possible with the same tools and this is where the results of this work come into play (Theorem 8, Corollaries 2 and 3). Second we deal with non-semimartingale models (Theorems 10 and 11): Fractional Brownian motion and Iterated Brownian motion.

### 4.1 Martingale at random times

On a filtered probability space which filtration satisfies the usual condition, let $\left(M_{t}\right)_{t \geq 0}$ be a $\mathbb{R}^{d}$-valued continuous martingale, which component-wise bracket is of the form $\left\langle M^{(i)}\right\rangle_{t}=\int_{0}^{t} m_{s}^{(i)} \mathrm{d}$ for a progressively measurable process $m^{(i)}$ bounded by a finite constant $C^{(M)}$.

Proposition 6. The measurable mapping $(\omega, t) \mapsto M_{t}(\omega)$ satisfies $(\mathbf{H} 1)$ and (H2) for any $\kappa \in\left(0, \frac{1}{2}\right)$.

Proof. Assumption (H1) is easily checked with $\alpha_{p}^{(\mathbf{H 1 )}}=1 / 2$ owing to Burkholder-Davis-Gundy inequalities [RY99, Theorem 4.1, p.160]. Let us now turn to the verification of (H2) for any $\kappa \in\left(0, \frac{1}{2}\right)$; let $\kappa$ be such a parameter. First, observe that it is enough to prove the $\mathbf{L}_{p}$-estimates in (H2) for $p$ large enough since they are automatically satisfied for smaller $p$ with the same constants, using the immediate inequality $\|Z\|_{\mathbf{L}_{q}} \leq\|Z\|_{\mathbf{L}_{p}}$ for $q \leq p$. Therefore, we now consider $p$ large enough such that $\kappa<1 / 2-1 / p$ (in particular $p>2$ ). A direct application of Burkholder-Davis-Gundy inequalities gives

$$
\left\|M_{t}-M_{s}\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(M)}|t-s|^{1 / 2}, \quad \forall s, t \in \mathbb{R}^{+}
$$

for some finite constant $C_{p}^{(M)}$ depending on $d, C_{p}^{\mathrm{BDG}}$ and $C^{(M)}$. Thus, (6) is fulfilled with $\beta^{(G)}=1 / 2$ and $\tau^{(G)}=0$ : from Theorem 2 we deduce (7) with $\beta=\kappa \in$ $\left(0, \beta^{(G)}-1 / p\right)$ : thus, (H2) holds for $F=M$ with $\alpha_{p}^{(\mathbf{H 2})}=1 / 2-\kappa$. We are done.

Theorem 8. Let $\theta^{N}$ and $\theta$ be random times with finite moments at any order, uniformly bounded w.r.t. $N$. For any $p>0$, any $\kappa \in(0,1 / 2)$ and $q>\kappa p$, there is a constant $c_{p, \kappa, q}$ such that

$$
\begin{equation*}
\left\|M_{\theta^{N}}-M_{\theta}\right\|_{\mathbf{L}_{p}} \leq c_{p, \kappa, q}\left\|\theta^{N}-\theta\right\|_{\mathbf{L}_{q}}^{\kappa} \tag{35}
\end{equation*}
$$

Proof. Let $\kappa \in\left(0, \frac{1}{2}\right)$. In view of Proposition 6 , we can apply Theorem 1 which reads in the current setting $\left(F=F^{N}=M\right.$ and $C_{p}^{N,(\mathbf{H 3})}=0, \Theta=\theta$ and $\left.\Theta^{N}=\theta^{N}\right)$

$$
\left\|M\left(\theta^{N}\right)-M(\theta)\right\|_{\mathbf{L}_{p}} \leq c_{(3)}\left\|\theta^{N}-\theta\right\|_{\mathbf{L}_{\kappa p_{2}}}^{\kappa}
$$

for any parameter $p_{2}>p$. This allows the choice $p_{2}=q / \kappa$ which leads to the advertised estimate.

As a comparison, we state a similar result available when $\theta^{N}$ and $\theta$ are stopping times. The proof is based on the Burkholder-Davis-Gundy inequalities applied to the martingale $N_{t}=M_{\theta^{N} \vee \theta \wedge t}-M_{\theta^{N} \wedge \theta \wedge t}$, we leave details to the reader.

Proposition 7. Let $\theta^{N}$ and $\theta$ be stopping times with finite moments at any order, uniformly bounded w.r.t. $N$. For any $p>0$, there is a constant $c_{p}$ such that

$$
\left\|M_{\theta^{N}}-M_{\theta}\right\|_{\mathbf{L}_{p}} \leq c_{p}\left\|\theta^{N}-\theta\right\|_{\mathbf{L}_{p / 2}}^{1 / 2} .
$$

Observe that the exponent of the $L_{p}$-norms of $\theta^{N}-\theta$ is slightly better in Proposition 7 than in Theorem 8 but the scope of applicability is narrower because of the restriction to stopping times in Proposition 7.

With a result like Theorem 8 at hand, we can study quite efficiently some non trivial approximation problems. Consider the approximation of the maximum of a scalar continuous martingale $M(d=1)$ on the time interval $[0, T]$ (with $0<T<$ $+\infty)$ upon discrete time monitoring. Set $\tau^{*}:=\inf \left\{t \in[0, T]: M_{t}=\max _{s \leq T} M_{s}\right\}$ for the first time at which $M$ reaches its maximum on $[0, T]$; clearly this is not a stopping time for the underlying filtration. It may happen that the maximum is
achieved several times (although it is a.s. unique for the Brownian motion), which justifies why we choose the first time.

Generally speaking, computing exactly $\tau^{*}$ is challenging: in practice, it can be approximated on a grid $\left(t_{i}:=i T / N\right)_{0 \leq i \leq N}($ with $N \geq 1)$ by the discrete time $\tau^{*, N}:=\inf \left\{t_{i} \in[0, T]: M_{t_{i}}=\max _{t_{j} \leq T} M_{t_{j}}\right\}$. In the Brownian case for $M$, we know that the error $M_{\tau^{*}}-M_{\tau^{*}, N}$ converges to 0 at rate $\sqrt{N}$, see [AGP95, Theorem 1, Lemma 6] for details. Owing to Theorem 8, we can prove that in the current more general case the strong error is of order $N^{-\kappa}$ for any $\kappa<1 / 2$.

Corollary 2. For any $\kappa \in(0,1 / 2)$ and any $p>0$, we have

$$
\left\|M_{\tau^{*}}-M_{\tau^{*}, N}\right\|_{\mathbf{L}_{p}}=O\left(N^{-\kappa}\right) .
$$

Proof. Set $\theta:=\tau^{*}$ and define $\theta^{N}$ as the closest point to $\theta$ on the discrete grid. Observe that it may be different from $\tau^{*, N}$, but anyhow we have

$$
\begin{aligned}
0 \leq M_{\tau^{*}}-M_{\tau^{*}, N} & \leq M_{\theta}-M_{\theta^{N}} \\
\left|\theta-\theta^{N}\right| & \leq T / N .
\end{aligned}
$$

The proof is finished in view of (35).

### 4.2 Local times at random time and random level

In this paragraph, Theorem 1 is applied to the case where the random map $F($. is the local time $\{L(t, x) ; x \in \mathbb{R}, t \geq 0\}$ of a scalar Brownian motion $W$, and where $\Theta=(\tau, \xi)$ is a random pair (time, level). Recall that $L$ is defined by the occupation-time formula

$$
\int_{\mathbb{R}} f(x) L(t, x) \mathrm{d} x=\int_{0}^{t} f\left(W_{s}\right) \mathrm{d} s
$$

for any $t \geq 0$ and any measurable function $f: \mathbb{R} \mapsto \mathbb{R}^{+}$. By [RY99, Theorem 1.7 p.225], $L$ has a bi-continuous modification that we consider from now.

Approximating Brownian local times at deterministic or random point $\Theta=(\tau, \xi)$ is interesting on its own and it has nice applications: for instance, we refer to the Ray-Knight theorems [RY99, Chapter XI], where local times at some random time $\tau$ are related to Bessel processes. A second example is the toy model of [FP11] where $W$ and $W+\alpha L(., a)$ respectively model the value of a stock in absence or presence of a proportion $\alpha$ of investors buying as soon as the price falls below $a$.

There exists several approximation schemes for the Brownian local time, see [Kho94] and references therein; in [Kho94], using the number of up-crossings to approximate the local time, sharp almost sure convergence rates in sup-norm are established. We prefer to take advantage of the recent work [OS14], deriving $\mathbf{L}_{p^{-}}$ estimates which fit well our setting. Let us recall their result by following closely their presentation. For a fixed positive integer $N$, we define $T_{0}^{N}:=0$ a.s. and

$$
T_{i}^{N}:=\inf \left\{t>T_{i-1}^{N}:\left|W_{t}-W_{T_{i-1}^{N}}\right|=2^{-N}\right\}, \quad i \geq 1 .
$$

Let $W^{N}$ denote the symmetric random walk (with non-equidistant jump times) defined by

$$
W^{N}(t):=\sum_{i=1}^{\infty} 2^{-N} \eta_{i}^{N} \mathbf{1}_{\left\{T_{i}^{N} \leq t\right\}}, \quad t \geq 0,
$$

where $\eta_{i}^{N}:=\operatorname{Sign}\left(W_{T_{i}^{N}}-W_{T_{i-1}^{N}}\right)$, for $i \geq 1$. Now, for a given $x \in \mathbb{R}$, let $j_{N}(x)$ be the unique integer such that $\left(j_{N}(x)-1\right) 2^{-N}<x \leq j_{N}(x) 2^{-N}$ and define

$$
\begin{aligned}
u\left(j_{N}(x) 2^{-N}, N, t\right):= & \text { the number of up-crossings of } W^{N} \text { from } \\
& \left(j_{N}(x)-1\right) 2^{-N} \text { to } j_{N}(x) 2^{-N} \text { before time } t .
\end{aligned}
$$

Finally, for $(t, x) \in[0,+\infty) \times \mathbb{R}$ set

$$
L^{N}(t, x):=22^{-N} u\left(j_{N}(x) 2^{-N}, N, t\right) .
$$

This gives the candidate for strongly approximating $L$, this is the next statement.
Theorem 9 ([OS14, Theorem 2.2]). For any fixed $T>0$ and any $p>0$, we have

$$
\sup _{N \geq 0}\left\|\sup _{0 \leq t \leq T} \sup _{x \in \mathbb{R}} \frac{\left|L^{N}(t, x)-L(t, x)\right|}{\left(2^{-N} \log \left(2^{N}\right)\right)^{\frac{1}{2}}}\right\|_{\mathbf{L}_{p}}<+\infty .
$$

In the following, we restrict to bounded time (say by $T$ ), or equivalently we consider $F^{N}(t, x)=L^{N}(t \wedge T, x)$ and $F(t, x)=L(t \wedge T, x)$ for any $t \geq 0$. Theorem 9 ensures that (H3) (with $\left.\alpha_{p}^{(\mathbf{H} 3)}=0\right)$ holds true for such $F^{N}$ and $F$.

We now investigate the validity of (H2). We start with a standard result.
Lemma 3. For any $p>0$, there is a finite constant $C_{p, T}$ such that

$$
\|L(t \wedge T, x)-L(s \wedge T, y)\|_{\mathbf{L}_{p}} \leq C_{p, T}\left(|t-s|^{\frac{1}{2}}+|x-y|^{\frac{1}{2}}\right)
$$

for any $x, y \in \mathbb{R}$ and any $t, s \in \mathbb{R}^{+}$.
The above estimation w.r.t. time follows easily from the Tanaka formula, the one w.r.t. space is stated in [RY99, Exercise 1.33 p.238]. We mention that similar controls in the more general case of continuous local martingales are proved in [BY82]. Then, as a consequence of Theorem $2,(t, x) \mapsto L(t \wedge T, x)$ is locally $\kappa$-Hölder, for any $\kappa \in] 0, \frac{1}{2}\left[\right.$, and (H2) holds for such $\kappa$ and $\alpha_{p}^{(\mathbf{H} 2)}=\frac{1}{2}-\kappa$. Similarly Corollary 1 implies (H1). Therefore, we can apply Theorem 1 and it gives the following result.

Corollary 3. Let $T>0$ be fixed. Let $\left(\tau^{N}, \tau\right)$ be finite random times and let $\left(\xi^{N}, \xi\right)$ be two scalar random variables with finite $\mathbf{L}_{p}$-norms (for any $p>0$ and uniformly in $N)$. Then for any $p>0, p_{2}>p$ and $\kappa \in(0,1 / 2)$, there is a finite constant $c$ such that

$$
\begin{aligned}
& \left\|L^{N}\left(\tau^{N} \wedge T, \xi^{N}\right)-L(\tau \wedge T, \xi)\right\|_{\mathbf{L}_{p}} \\
& \quad \leq c\left[2^{-\frac{N}{2}}\left(\log \left(2^{N}\right)\right)^{\frac{1}{2}}+\left\|\tau^{N} \wedge T-\tau \wedge T\right\|_{\mathbf{L}_{\kappa p_{2}}}^{\kappa}+\left\|\xi^{N}-\xi\right\|_{\mathbf{L}_{\kappa p_{2}}}^{\kappa}\right] .
\end{aligned}
$$

### 4.3 Fractional Brownian motion at random times

Models based on fractional fields are now popular in physics, natural sciences, economy and finance among other fields, see [CI13]. As an application of our approximation results, we consider the model of fractional Brownian motion (fBM in short), introduced in [MV68], which we denote by $\left(B_{t}^{(H)}\right)_{t \in \mathbb{R}}$. It is parametrized by its Hurst exponent $H \in(0,1)$, it is a $\mathbb{R}$-valued Gaussian process, centered with covariance function

$$
\mathbb{E}\left(B_{t}^{(H)} B_{s}^{(H)}\right)=\frac{1}{2}\left(|t|^{2 H}+|s|^{2 H}-|t-s|^{2 H}\right) .
$$

Remarkably, this is not a semimartingale, which contrasts with previous diffusion or martingale models we have considered so far. It is also used in financial modeling, see [CR98]. Our aim is to study the strong approximation of $B_{\Theta \wedge T}^{(H)}$ where $\Theta \geq 0$ is a random time and $T$ is fixed.

There are multiple possible approximations of $B^{(H)}$ (see [Sza01], [HMBL14] and references therein), we do not enter into details. For the sake of conciseness, assume directly that $B^{(H)}$ is approximated on $[0, T]$ by a scheme $B^{(H), N}$ depending on an algorithm parameter $N \rightarrow+\infty$. Assume the existence of a non-negative sequence $\left(\varepsilon_{N}\right)_{N \geq 1}$ converging to 0 such that for any $p \geq 1$,

$$
\left\|\sup _{0 \leq t \leq T}\left|B_{t}^{(H), N}-B_{t}^{(H)}\right|\right\|_{\mathbf{L}_{p}}=O\left(\varepsilon_{N}\right) .
$$

For instance, in [HMBL14] we have $\varepsilon_{N}=\frac{\sqrt{\log N}}{N^{H A(1-H)}}$, which readily follows from their Theorem 6.1 (with a restriction to rational numbers $t$ in the above sup). The scheme described in [Sza01] converges ${ }^{2}$ at rate $\varepsilon_{N}=\frac{\log N}{N^{H \wedge \frac{1}{2}}}$ : this can be easily derived from [Sza01, Theorem 3].

Now, remind that for any $p>0,\left\|B_{t}^{(H)}-B_{s}^{(H)}\right\|_{\mathbf{L}_{p}}=c_{p}|t-s|^{H}$ (for any $s, t \geq 0$ ) for some constant $c_{p}$ : therefore, Theorem 2 yields that $(\mathbf{H} 1)-(\mathbf{H} 2)$ are fulfilled for $t \mapsto B_{t}^{(H)}$ for any $\kappa<H$. If $\Theta$ is additionally approximated by $\Theta^{N}$, we obtain the global error estimates as a consequence of Theorem 1.

Theorem 10. With the above notation, for any $\kappa \in(0, H)$ and any $p>0$ we have

$$
\left\|B_{\Theta^{N} \wedge T}^{(H), N}-B_{\Theta \wedge T}^{(H)}\right\|_{\mathbf{L}_{p}} \leq O\left(\varepsilon_{N}+\left\|\Theta^{N} \wedge T-\Theta \wedge T\right\|_{\mathbf{L}_{p H}}^{\kappa}\right) .
$$

### 4.4 Diffusion process in Brownian time

During the two last decades, there has been an increasing interest for studying Diffusion processes in Brownian time. It dates back to the work by Funaki [Fun79], and it is furthermore studied in [Bur93] under the name Iterated Brownian Motion (IBM in short) as

$$
\begin{equation*}
Z_{t}=\tilde{B}_{B_{t}} \tag{36}
\end{equation*}
$$

where $\left(\tilde{B}_{t}\right)_{t \in \mathbb{R}}$ is a two-sided $\mathbb{R}^{d}$-valued Brownian motion and $\left(B_{t}\right)_{t \geq 0}$ is a scalar Brownian motion independent of $\tilde{B}$. It serves, for instance, for modeling the Brownian motion in a Brownian crack [BK98] (limit of a Brownian motion reflected in a

[^2]Wiener sausage of width shrinking to 0 ). This process has nice properties like selfsimilarity, stationary increments and $\alpha$-Hölder continuous paths (with $\alpha<1 / 4$ ), it is not a semimartingale. As variants of this model, let us mention the $n$-times iteration of Brownian motion, studied in [CK14] (with $n \rightarrow+\infty$ ), and the case where $\tilde{B}$ is replaced by a fractional Brownian motion with Hurst index $H \in(0,1)$ [NZ14]. Moreover, the IBM (36) allows to represent Feynman-Kac solutions of fourth-order PDEs (see [Fun79], [All02] and references therein): indeed, under appropriate conditions on $f, u(t, x)=\mathbb{E}\left(f\left(x+Z_{t}\right)\right)$ solves the Bi-Laplacian equation

$$
\left\{\begin{array}{l}
\partial_{t} u(t, x)=\frac{\Delta f(x)}{\sqrt{8 \pi t}}+\frac{1}{8} \Delta^{2} u(t, x), \quad t>0, x \in \mathbb{R}^{d}, \\
u(0, x)=f(x), \quad x \in \mathbb{R}^{d} .
\end{array}\right.
$$

Replacing $B$ in (36) by a diffusion process leads to more general fourth-order PDEs [AZ01]. To account even for greater generality, we now consider diffusion process in diffusion time and study its strong approximation. This writes:

$$
Z_{t}=X_{\left|Y_{t}\right|}
$$

where

$$
\begin{array}{r}
X_{t}=x+\int_{0}^{t} \mu\left(X_{s}\right) \mathrm{d} s+\int_{0}^{t} \sum_{i=1}^{q} \sigma_{i}\left(X_{s}\right) \mathrm{d} W_{s}^{i}, \\
Y_{t}=y+\int_{0}^{t} b\left(Y_{s}\right) \mathrm{d} s+\int_{0}^{t} \sum_{i=1}^{q} \gamma^{i}\left(Y_{s}\right) \mathrm{d} W_{s}^{i} .
\end{array}
$$

Here, $W=\left(W^{1}, \ldots, W^{q}\right)$ is $q$-dimensional Brownian motion, $X$ takes values in $\mathbb{R}^{d}$ and $Y$ in $\mathbb{R}^{d^{\prime}}$. Up to the $\mathbb{R}^{d^{\prime}}$-norm term $|$.$| which permits to avoid the use of two-$ sided process, observe that it includes the model (36) by choosing appropriately the coefficients $\mu, b, \sigma, \gamma$ and by setting $W=(\tilde{B}, B)$. Here the coefficients $\mu, b, \sigma, \gamma$ do not depend in time, this is only for the sake of simplification of the statement.

Now consider the continuous-time Euler scheme $X^{\delta}$ and $Y^{\delta}$, with time step $\delta \in(0,1]$, associated to $X$ and $Y$ : they are both defined similarly to (13) by setting $\tau_{t}:=i \delta$ for $i \delta \leq t<(i+1) \delta$. We define $Z_{t}^{\delta}:=X_{\left|Y_{t}^{\delta}\right|}^{\delta}$. Our main result states that the strong convergence order is almost $\frac{1}{4}$.

Theorem 11. Assume that $b, \mu, \sigma, \gamma$ are bounded Lipschitz functions. Let $\kappa \in\left(0, \frac{1}{4}\right)$ : then, for any $p>0$ and any $T>0$, we have

$$
\sup _{t \in[0, T]}\left\|X_{\left|Y_{t}^{\delta}\right|}^{\delta}-X_{\left|Y_{t}\right|}\right\|_{\mathbf{L}_{p}}=O\left(\delta^{\kappa}\right)
$$

Before proving the above result, observe that the effective simulation of $Z_{t}^{\delta}$ can be easily performed. First, sample the Brownian increments $\left(W_{j \delta}-W_{(j-1) \delta}\right)_{j \geq 1}$ up to the index $i:=\tau_{t} / \delta$, as well as $W_{t}-W_{i \delta}$ : this is sufficient to obtain the Euler scheme time $t^{\prime}:=\left|Y_{t}^{\delta}\right|$. Then to simulate $X_{t^{\prime}}^{\delta}$, two cases have to be considered. Set $i^{\prime}:=\tau_{t^{\prime}} / \delta$ and denote by $\mathcal{G}$ the sigma-field generated by the $(i+1)$ previous random variables.

1. If $t^{\prime} \geq t$, sample additionally $W_{(i+1) \delta}-W_{t}, W_{(i+2) \delta}-W_{(i+1) \delta}, \ldots, W_{i^{\prime} \delta}-$ $W_{\left(i^{\prime}-1\right) \delta}, W_{t^{\prime}}-W_{i^{\prime} \delta}$, which are (conditionally to $\mathcal{G}$ ) independent centered Gaussian random variables, with independent components having variance equal to the time-increments. Their simulation is thus straightforward and yields $X_{t^{\prime}}^{\delta}$.
2. If $t^{\prime}<t$, it is enough to simulate $W_{t^{\prime}}-W_{i^{\prime} \delta}$ conditionally on $\mathcal{G}$, i.e. the marginal distribution of a Brownian bridge. Namely, if $t^{\prime}<i \delta$, then it is Gaussian distributed, with mean $\left(W_{\left(i^{\prime}+1\right) \delta}-W_{i^{\prime} \delta}\right) \frac{\left(t^{\prime}-i^{\prime} \delta\right)}{\delta}$ and variance $\frac{\left.\left(t^{\prime}-i^{\prime} \delta\right)\left(i^{\prime}+1\right) \delta-t^{\prime}\right)}{\delta}$ (component by component). Otherwise, if i $\delta \leq t^{\prime}<t$, the mean and variance are adjusted to $\left(W_{t}-W_{i \delta}\right) \frac{\left(t^{\prime}-i \delta\right)}{(t-i \delta)}$ and $\frac{\left(t^{\prime}-i \delta\right)\left(t-t^{\prime}\right)}{t-i \delta}$. The simulation of $X_{t^{\prime}}^{\delta}$ readily follows.

Proof. It is enough to prove the result for $p \geq 2$. One knows [BGG14, Lemma A.2] that there exists a constant $c_{p}^{(X)}>0$ (depending on $d, q$, on the bounds of $\mu, \sigma$ and their Lipschitz constants) such that for any $T>0$

$$
\begin{equation*}
\left\|\sup _{t \in[0, T]}\left|X_{t}^{\delta}-X_{t}\right|\right\|_{\mathbf{L}_{p}} \leq c_{p}^{(X)} e^{c_{P}^{(X)}} \delta^{\frac{1}{2}} \tag{37}
\end{equation*}
$$

As the reader may guess, the exponential term w.r.t. time comes from an application of Gronwall's lemma in the derivation of error estimates. Note that this exponential factor does not enable us to deduce that $\left.\left(X_{s}^{\delta}-X_{s}\right)\right|_{s=\left|Y_{t}\right|}$ belongs to any $\mathbf{L}_{p}$, with good $\mathbf{L}_{p}$-estimates, by applying Proposition 1 since the latter requires a polynomial growth in the stochastic argument (here $T$ ). But actually, we can modify Proposition 1 by taking advantage of the finite exponential moments of $\left|Y_{t}\right|$ (because $b$ and $\gamma$ are bounded). This is the next statement which we will prove at the end.
Proposition 8. Let $E$ be an Euclidean space and $G$ be a $\mathcal{F} \otimes \mathcal{B}(E)$-measurable mapping taking values in $\mathcal{E}$ such that, for any $p>0$, there exist constants $\alpha_{p}^{(G)} \in$ $[0,+\infty)$ and $C_{p}^{(G)} \in[0,+\infty)$ for which

$$
\left\|\sup _{|x| \leq \lambda}|G(\cdot, x)|\right\|_{\mathbf{L}_{p}} \leq C_{p}^{(G)} \exp \left(\alpha_{p}^{(G)} \lambda\right), \quad \forall \lambda \geq 0
$$

Then, for any $E$-valued random variable $\xi$ with exponential moments, $\omega \mapsto G(\omega, \xi(\omega)) \in$ $\mathbf{L}_{p}$ for any $p$, and for any $\rho>0$ and any finite conjugate exponents $r$ and $s$, we have

$$
\|G(\xi)\|_{\mathbf{L}_{p}} \leq C_{p r}^{(G)}\left(e^{\rho r}-1\right)^{-1 /(p r)}\left(\mathbb{E}\left(e^{\left(p \alpha_{p r}^{(G)}+\rho\right) s|\xi|}\right)\right)^{1 /(p s)}
$$

Combine (37) with the above (for $\rho=1, r=s=2$ ): this proves that

$$
\left\|X_{\left|Y_{t}\right|}^{\delta}-X_{\left|Y_{t}\right|}\right\|_{\mathbf{L}_{p}} \leq c_{2 p}^{(X)} \delta^{\frac{1}{2}}\left(e^{2}-1\right)^{-1 /(2 p)}\left(\mathbb{E}\left(e^{2\left(p c_{2 p}^{(X)}+1\right)\left|Y_{t}\right|}\right)\right)^{1 /(2 p)}
$$

Since exponential moments of $Y_{t}$ are bounded locally uniformly in time, we get

$$
\begin{equation*}
\sup _{t \in[0, T]}\left\|X_{\left|Y_{t}\right|}^{\delta}-X_{\left|Y_{t}\right|}\right\|_{\mathbf{L}_{p}}=O\left(\delta^{\frac{1}{2}}\right) \tag{38}
\end{equation*}
$$

We now handle the difference $X_{\left|Y_{t}^{\delta}\right|}^{\delta}-X_{\left|Y_{t}\right|}^{\delta}$. As in the proof of Proposition 6 (since $\mu$ and $\sigma$ are bounded), we can easily prove that $\left\|X_{u}^{\delta}-X_{v}^{\delta}\right\|_{\mathbf{L}_{2 p}} \leq C_{2 p}^{(X)}|u-v|^{\frac{1}{2}}(1+$
$|u|+|v|)^{\frac{1}{2}}$ for any $u, v \geq 0$ for a constant $C_{2 p}^{(X)}$ depending only on $p, \mu$ and $\sigma$. Let $\beta \in(0,1 / 2-1 /(2 p))$ : then an application of Theorem 2 gives the existence of another constant $\tilde{C}_{2 p}^{(X)}$ such that

$$
H_{\beta}(\lambda):=\sup _{u \neq v, 0 \leq u, v \leq \lambda} \frac{\left|X_{u}^{\delta}-X_{v}^{\delta}\right|}{|u-v|^{\beta}}
$$

satisfies

$$
\left\|H_{\beta}(\lambda)\right\|_{\mathbf{L}_{2 p}} \leq \tilde{C}_{2 p}^{(X)} \lambda^{1-\beta}, \quad \forall \lambda \geq 1
$$

Now, note that $\sup _{\delta \in(0,1], t \in[0, T]}\left\|\left|Y_{t}^{\delta}\right| \vee\left|Y_{t}\right|\right\|_{\mathbf{L}_{q}}<+\infty$ for any $q>0$ (see (15) and (24) for $Y$ ), so that Proposition 1 yields

$$
\sup _{\delta \in(0,1], t \in[0, T]}\left\|H_{\beta}\left(\left|Y_{t}^{\delta}\right| \vee\left|Y_{t}\right|\right)\right\|_{\mathbf{L}_{3 p / 2}}<+\infty
$$

As a consequence, and using (37) (available for $Y$ and any $p$ ), we deduce

$$
\left\|X_{\left|Y_{t}^{\delta}\right|}^{\delta}-X_{\left|Y_{t}\right|}^{\delta}\right\|_{\mathbf{L}_{p}} \leq\left\|\left|\left|Y_{t}^{\delta}\right|-\left|Y_{t}\right|\right|^{\beta}\right\|_{\mathbf{L}_{3 p}}\left\|H_{\beta}\left(\left|Y_{t}^{\delta}\right| \vee\left|Y_{t}\right|\right)\right\|_{\mathbf{L}_{3 p / 2}} \leq c \delta^{\beta / 2}
$$

where the constant $c$ is uniform in $t \in[0, T]$ and $\delta \in(0,1]$. Gathering the above with (38), we obtain

$$
\sup _{t \in[0, T]}\left\|X_{\left|Y_{t}^{\delta}\right|}^{\delta}-X_{\left|Y_{t}\right|}\right\|_{\mathbf{L}_{p}}=O\left(\delta^{\beta / 2}\right) .
$$

So far, $\beta$ is in $(0,1 / 2-1 /(2 p))$ : hence, it is true for any $\beta<1 / 2$ provided that $p$ is large enough. It remains valid for smaller $p$ by stability of $\mathbf{L}_{p}$-norms. Theorem 11 is proved.

Proof of Proposition 8. We adjust the proof of Proposition 1 to the current exponential controls. Briefly, we have

$$
\begin{aligned}
\mathbb{E}\left(|G(\cdot, \xi)|^{p}\right) & \leq\left[C_{p r}^{(G)}\right]^{p} \sum_{n \geq 1} e^{-\rho n} e^{p \alpha_{p r}^{(G)} n} e^{\rho n} \mathbb{P}(n-1 \leq|\xi|<n)^{1 / s} \\
& \leq\left[C_{p r}^{(G)}\right]^{p}\left(\sum_{n \geq 1} e^{-\rho r n}\right)^{1 / r}\left(\sum_{n \geq 1} e^{\left(p \alpha_{p r}^{(G)}+\rho\right) s n} \mathbb{P}(n-1 \leq|\xi|<n)\right)^{1 / s} \\
& \leq\left[C_{p r}^{(G)}\right]^{p}\left(\frac{e^{-\rho r}}{1-e^{-\rho r}}\right)^{1 / r}\left(\mathbb{E}\left(e^{\left(\alpha_{p r}^{(G)}+\rho\right) s|\xi|}\right)\right)^{1 / s}
\end{aligned}
$$

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[^1]:    ${ }^{1}$ That is for any compact $K$ of $\mathbb{R}^{d}$ there exists a finite positive random variable $C(K)$ such that for any $x, y \in K$ we have $\left|\nabla X_{t}(x, \omega)-\nabla X_{t}(y, \omega)\right| \leq C(K, \omega)|x-y|^{\varepsilon}$ a.s., see [Kun97, Chapters 3 and 4] for details.

[^2]:    ${ }^{2}$ to be precise, their scheme converges towards some fBM and not necessarily $B^{(H)}$.

