Approximation of the Invariant Measure with an Euler Scheme for Stochastic PDEs Driven by Space-Time White Noise

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Abstract In this article, we consider a stochastic PDE of parabolic type, driven by a space-time white-noise, and its numerical discretization in time with a semi-implicit Euler scheme. When the nonlinearity is assumed to be bounded, then a dissipativity assumption is satisfied, which ensures that the SDPE admits a unique invariant probability measure, which is ergodic and strongly mixing—with exponential convergence to equilibrium. Considering test functions of class C^2 , bounded and with bounded derivatives, we prove that we can approximate this invariant measure using the numerical scheme, with order 1/2 with respect to the time step.

Keywords Stochastic partial differential equations • Invariant measures and Ergodicity • Weak approximation • Euler scheme

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1 Introduction

In this paper, we are interested in the discretization in time of the following stochastic reaction-diffusion equation

$$\frac{\partial y(t,\xi)}{\partial t} = \frac{\partial^2 y(t,\xi)}{\partial \xi^2} + g(\xi, y(t,\xi)) + \frac{\partial \omega(t,\xi)}{\partial t},\tag{1}$$

for $t \ge 0, \xi \in (0, 1)$, with the initial condition $y(0, \xi) = y(\xi)$, and homogeneous Dirichlet boundary conditions y(t, 0) = y(t, 1) = 0. The stochastic perturbation $\frac{\partial \omega(t,\xi)}{\partial t}$ is a space-time white noise: the rigorous interpretation of Eq. 1 is given by an abstract evolution (2)—in the sense of [3]—in the Hilbert space $H = L^2(0, 1)$,

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ENS Cachan Bretagne - IRMAR, Université Rennes 1, Avenue Robert Schuman, 35170 Bruz, France e-mail: charles-edouard.brehier@bretagne.ens-cachan.fr driven by a Wiener cylindrical process—see Section 2.3. We assume that g is of class C^2 with respect to y and measurable with respect to ξ . Moreover we assume uniform bounds on g and on its derivatives, uniformly with respect to ξ ; more details are given in Example 2.7.

The numerical approximation of stochastic equations has been extensively studied during the last thirty years. First we recall that one can look at strong approximation results—when the trajectories of the continuous and the discrete-time processes are compared—or at weak approximation results—when the laws at a fixed time are compared. The simplest method in the case of SDEs is the Euler-Maruyama scheme; it is built as a straightforward extension of the well-known explicit Euler method for ODEs, which is of order 1: if we consider in \mathbb{R}^d a SDE with regular coefficients

$$dX_t = f(X_t)dt + \sigma(X_t)dB_t, X_0 = x,$$

its numerical approximation is defined for a given step-size Δt by

$$X_0 = x,$$

$$X_{n+1} = X_n + f(X_n)\Delta t + \sigma(X_n)(B_{t_{n+1}} - B_{t_n})$$

Due to the regularity properties of the Brownian Motion, this scheme is in general only of order 1/2 in the strong sense—i.e. for $n\Delta t \leq T$ we have $\mathbb{E}|X_{n\Delta t} - X_n|^2 \leq C(T)\Delta t$ —while it is of order 1 in the weak sense—i.e. for test functions ϕ of class C^3 , with bounded derivatives, we have a bound $|\mathbb{E}\phi(X_{n\Delta t}) - \mathbb{E}\phi(X_n)| \leq C(T, \phi)\Delta t$. The idea for proving that weak order is 1 and not 1/2 is to consider the Kolmogorov equation associated with the process X_t , which is satified by the function $(t, x) \mapsto \mathbb{E}\phi(X(t, x))$ —see [28, 30]. The books [14] and [22]—see also [21]—contain various numerical schemes—like the well-known Milstein scheme, some implicit schemes, and methods based on stochastic Taylor expansions—with their order of convergence in both strong and weak senses.

Numerical methods for SPDEs like (Eq. 1) need discretization both in time and in space. For example, time discretization leads to explicit or implicit methods, while discretization in space can be done with finite difference or finite element methods. Basically, the result for space-time white noise driven equations is convergence with strong order 1/2 in space and only 1/4 in time—under some Courant-Friedrichs-Lewy conditions when necessary: see [5, 10–13, 31].

If we look at the abstract formulation of Eq. 1 in the Hilbert space H, results have also been proved for the time discretization using semi-implicit Euler schemes: the strong order of convergence is 1/4—see [25]—while the weak order of convergence is 1/2—see [6]. Moreover in [8], the authors have studied weak convergence in the case of linear equations when using a finite element method in space. We follow here the framework of [6]: we consider the stochastic evolution equation in H and we use a time discretization with a semi-implicit Euler scheme, with no discretization in space.

In this work, we are interested in the behaviour of the weak convergence estimates when the final time T goes to infinity: can we replace constants $C(T, \phi)$ by a constant $C(\phi)$ independent from T? Passing to the limit, we thus ask the more general following question: can we use a numerical scheme to approximate the invariant probability measure of the continuous time process—which is assumed to be unique, ergodic and with exponential convergence to equilibrium? The SDE case has been studied with various methods: in [29], the weak error analysis is made by showing that the time derivatives of the solution of the Kolmogorov equation are exponentially decreasing in time; in [17], some general conditions are given for the ergodicity of the numerical scheme, thanks to the theory of geometric ergodicity of Markov Chains. In [19], the approximation result is shown thanks to the use of a Poisson equation.

In the case of SPDEs, general ergodicity results have only been obtained recently—see Section 4.3—and the problem of approximation of invariant measures by numerical schemes has not been studied yet.

We prove the following result:

Theorem 1.1 For any $0 < \kappa < 1/2$, $\tau_0 > 0$ and for any C_b^2 function ϕ , there exists a constant C > 0 such that for any $m \ge 2$, $y \in H$ and $0 < \tau \le \tau_0$

$$|\mathbb{E}[\phi(Y(m\tau, y))] - \mathbb{E}[\phi(Y_m(\tau, y))]| \le C(1+|y|^3)(((m-1)\tau)^{-1/2+\kappa}+1)\tau^{1/2-\kappa}$$

The continuous process *Y* is defined by Eq. 2 below, and the numerical approximation $(Y_m(\tau, y))$ with time step τ and initial condition *y* is defined by Eq. 8.

We notice that the right-hand side of the previous estimate contains a singularity when m = 1, which is due to a lack of regularity of the infinite-dimensional processes—details are given in Section 6.1 below. However, if we look at the error at a fixed time $T = m\tau$, if τ is small enough we just need to change the constant C = C(T); moreover since we are interested in the behaviour when m goes to infinity, this term plays no role.

With the assumptions precised below, we show in Section 4.3 that the SPDE admits a unique invariant probability measure $\overline{\mu}$, which is ergodic and strongly mixing, with exponential convergence to equilibrium; nevertheless we can in general only show the existence, not the uniqueness, of invariant measures for the numerical approximation, and we can prove the following result:

Corollary 1.2 For any $0 < \kappa < 1/2$, $\tau_0 > 0$ and for any C_b^2 function ϕ , there exists constants c > 0, C > 0 such that for any $0 < \tau \le \tau_0$, any initial condition $y \in H$ and any $m \ge 1$

$$|\mathbb{E}[\phi(Y_m(\tau, y))] - \int_H \phi d\overline{\mu}| \le C(1 + |y|^3) \left(\frac{1}{m^{1/2-\kappa}} + \tau^{1/2-\kappa}\right) + C(1 + |y|^2)e^{-cm\tau}$$

Moreover, if μ^{τ} is an ergodic invariant probability measure of the numerical scheme $(Y_m(\tau, .))_{m \in \mathbb{N}}$, we have

$$\left|\int_{H}\phi d\mu^{\tau} - \int_{H}\phi d\overline{\mu}\right| \leq C\tau^{1/2-\kappa};$$

all the ergodic invariant measures of the numerical approximation are then close to the unique invariant probability measure of the continuous process.

Up to our knowledge, this is the first result of this kind for SPDEs.

The key point, like in [29], for obtaining bounds independent from the time $T = m\tau$, is to prove that derivatives of the solution *u* of the underlying Kolomogorov equation—mentioned above—decrease exponentially in time: this is done in Section 5.2, using a coupling method.

Moreover, an essential tool in [6] is the use of Malliavin calculus and of an integration by parts formula in order to transform a stochastic integral—which is

not regular enough in space in the infinite dimensional setting—into an expression which can be controlled; in Lemma 4.5, we see that the involved Malliavin derivatives may increase exponentially fast with respect to time. The solution is to separate the lack of regularity problem and this badly controlled growth with a decomposition of the interval into two parts: in the first one we need the integration by parts formula, and we can control the Malliavin derivatives, while in the other one we can directly give an appropriate bound. We also provide an improvement with respect to [6]: we can consider a more general nonlinear coefficient G—like a Nemytskii operator, see Example 2.7.

We also notice that some more general equations, with additive noise which is white in time but colored in space, can be studied with our method, with suitable assumptions on the coefficients. For a very smooth noise, the numerical analysis on finite time is easier to treat, but then ergodic properties and long-time behaviour require different techniques; this will be treated elsewhere.

The case of some equations with multiplicative noise which satisfy the Strong Feller Property is also covered by our technique of proof, but with some additional difficulties—see for instance the restrictive condition on the diffusion coefficient in [6]. Since all the necessary ideas are contained here and in [6], we only focus on the additive noise case.

The paper is organized as follows: in Section 2, we precise the assumptions made on the coefficients of the equations, and we define the numerical method in Section 3. In Section 4, we give some bounds on the solutions of the continuous and discrete equations, and we study their asymptotic behaviour: existence of invariant measures, and uniqueness for the continuous equation, following from Proposition 4.7. In Section 5, we explain the proof of Theorem 1.1, and we give a proof of Corollary 1.2; more precisely, in Section 5.2 we prove the exponential decreasing in time of the derivatives of the function u. Eventually Section 6 contains the proofs of the remaining estimates

2 Notations and Assumptions

Let *H* be a separable Hilbert space, with norm denoted by $|.|_H$ or simply |.|. We consider equations of the form

$$dY(t, y) = (BY(t, y) + G(Y(t, y)))dt + dW(t)$$

Y(0, y) = y. (2)

In the next paragraphs, we explain the assumptions on the linear operator B and on the nonlinear coefficient G; we also recall how the cylindrical Wiener process W is defined, and how we can construct solutions of this equation.

2.1 Test Functions

To quantify the weak approximation, we use test functions ϕ in the space $C_b^2(H, \mathbb{R})$ of functions from H to \mathbb{R} that are twice continuously differentiable, bounded, with first and second order bounded derivatives.

Remark 2.1 In the sequel, we often identify the first derivative $D\phi(x) \in \mathcal{L}(H, \mathbb{R})$ with the gradient in the Hilbert space *H*, and the second derivative $D^2\phi(x)$ with a linear operator on *H*, via the formulas:

$$< D\phi(x), h > = D\phi(x).h \text{ for every } h \in H$$
$$< D^2\phi(x).h, k > = D^2\phi(x).(h, k) \text{ for every } h, k \in H$$

In the sequel, we use the following notations:

$$\|\Phi\|_{\infty} = \sup_{x \in H} |\Phi(x)|_{H}$$
$$\|\Phi\|_{1} = \sup_{x \in H} |D\Phi(x)|_{H}$$
$$\|\Phi\|_{2} = \sup_{x \in H} |D^{2}\Phi(x)|_{\mathcal{L}(H)}$$

2.2 Assumptions on the Coefficients

2.2.1 The Linear Operator

We denote by $\mathbb{N} = \{0, 1, 2, ...\}$ the set of nonnegative integers. We suppose that the following properties are satisfied:

Assumptions 2.2

1. We assume that there exists a complete orthonormal system of elements of H denoted by $(f_k)_{k \in \mathbb{N}}$, and a non-decreasing sequence of real positive numbers $(\mu_k)_{k \in \mathbb{N}}$ such that:

$$Bf_k = -\mu_k f_k$$
 for all $k \in \mathbb{N}$.

2. The sequence (μ_k) goes to $+\infty$ and

$$\sum_{k=0}^{+\infty} \frac{1}{\mu_k^{\alpha}} < +\infty \Leftrightarrow \alpha > 1/2.$$

The smallest eigenvalue of -B is then μ_0 .

Example 2.3 The Eq. 1 enters in this framework: we can choose $B = \frac{D^2}{dx^2}$, with the domain $H^2(0, 1) \cap H^1_0(0, 1) \subset L^2(0, 1)$ —corresponding to homogeneous Dirichlet boundary conditions. In this case for any $k \in \mathbb{N}$ $\mu_k = \pi^2(k+1)^2$, and $f_k(\xi) = \sqrt{2} \sin((k+1)\pi\xi)$ —see [2].

For a $N \in \{1, 2, ...\}$, we denote by H_N the subspace of H spanned by $f_0, ..., f_{N-1}$, and by P_N the orthogonal projection of H onto H_N .

The domain D(B) of B is equal to $D(B) = \{y = \sum_{k=0}^{+\infty} y_k f_k \in H, \sum_{k=0}^{+\infty} (\mu_k)^2 |y_k|^2 < +\infty\}$. We can more generally define fractional powers of -B, for $b \in [0, 1]$:

$$(-B)^b y = \sum_{k=0}^{\infty} \mu_k^b y_k f_k \in H$$

with the domains

$$D(-B)^{b} = \left\{ y = \sum_{k=0}^{+\infty} y_{k} f_{k} \in H, \quad |y|_{b}^{2} := \sum_{k=0}^{+\infty} (\mu_{k})^{2b} |y_{k}|^{2} < +\infty \right\}.$$

When $b \in [0, 1]$, we can also define the spaces $D(-B)^{-b}$ and operators $(-B)^{-b}$, with norm denoted by $|.|_{-b}$; when $y = \sum_{k=0}^{+\infty} y_k f_k \in H$, we have $(-B)^{-b} y = \sum_{k=0}^{+\infty} \mu_k^{-b} y_k f_k$ and $|y|_{-b}^2 := \sum_{k=0}^{+\infty} (\mu_k)^{-2b} |y_k|^2$.

The semi-group $(e^{tB})_{t\geq 0}$ can be defined by the Hille-Yosida Theorem—see [2]. We use the following spectral formula: if $y = \sum_{k=0}^{+\infty} y_k f_k \in H$, then for any $t \geq 0$

$$e^{tB}y = \sum_{k=0}^{+\infty} e^{-\mu_k t} y_k f_k.$$

For any $t \ge 0$, e^{tB} is a continuous linear operator in H, with operator norm $e^{-\mu_0 t}$. The semi-group (e^{tB}) is used to define the solution $Z(t) = e^{tB}z$ of the linear Cauchy problem

$$\frac{dZ(t)}{dt} = BZ(t) \quad \text{with} \quad Z(0) = z.$$

To define solutions of more general PDEs of parabolic type, we use mild formulation, and Duhamel principle.

This semi-group enjoys some smoothing properties that we often use in this work. Basically we need the following properties, which are easily proved using the above spectral properties.

Proposition 2.4 Under Assumption 2.2, for any $\sigma \in [0, 1]$, there exists $C_{\sigma} > 0$ such that we have:

(1) for any t > 0 and $y \in H$,

$$|e^{tB}y|_{\sigma} \leq C_{\sigma}t^{-\sigma}e^{-\frac{\mu_0}{2}t}|y|_{H}.$$

(2) for any 0 < s < t and $y \in H$,

$$|e^{tB}y - e^{sB}y|_H \le C_{\sigma} \frac{(t-s)^{\sigma}}{s^{\sigma}} e^{-\frac{\mu_0}{2}s} |y|_H.$$

(3) for any 0 < s < t and $y \in D(-A)^{\sigma}$,

$$|e^{tB}y - e^{sB}y|_H \le C_{\sigma}(t-s)^{\sigma}e^{-\frac{\mu_0}{2}s}|y|_{\sigma}.$$

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2.2.2 The Nonlinear Operator

The nonlinear operator G is assumed to satisfy some general assumptions. In Example 2.7, we give the two main kind of operators that can be used in our framework.

Assumptions 2.5 The function $G : H \to H$ is assumed to be bounded and Lipschitz continuous. We denote by L_G the Lipschitz constant of G

We also define for each $N \ge 1$ a function $G_N : H_N \to H_N$, with $G_N(y) = P_N G(y)$ for any $y \in H_N$. We assume that each G_N is twice differentiable, and that we have the following bounds on the derivatives, uniformly with respect to N:

• There exists a constant C_1 such that for any $N \ge 1$, $y \in H_N$ and $h \in H_N$

$$|DG_N(y).h|_H \le C_1|h|_H.$$

• There exists $\eta \in [0, 1)$ and a constant C_2 such that for any $N \ge 1$, $y \in H_N$ and any $h, k \in H_N$ we have

$$|(-B)^{-\eta}D^2G_N(y).(h,k)| \le C_2|h|_H|k|_H.$$

• Moreover, there exists a constant C_3 such that for any $N \ge 1$, $y \in H_N$ and any $h, k \in H_N$

$$|D^2 G_N(y).(h,k)| \le C_3 |h|_{(-B)^{\eta}} |k|_H.$$

Since G is bounded, the following property is easily satisfied:

Proposition 2.6 (Dissipativity) *There exist* c > 0 *and* C > 0 *such that for any* $y \in D(B)$

$$\langle By + G(y), y \rangle \leq -c|y|^2 + C.$$
 (3)

We remark that we have uniform control with respect to the dimension N of the bounds on G_N and on its derivatives, and that Eq. 3 is also satisfied for G_N , with constants c and C independent from N.

Example 2.7 We give some fundamental examples of nonlinearities for which the previous assumptions are satisfied:

- A function $G: H \to H$ of class C^2 , bounded and with bounded derivatives, fits in the framework, with the choice $\eta = 0$.
- The function *G* can be a **Nemytskii** operator: let $g: (0, 1) \times \mathbb{R} \to \mathbb{R}$ be a measurable, bounded, function such that for almost every $\xi \in (0, 1)$ $g(\xi, .)$ is twice continuously differentiable, with uniformly bounded derivatives. Then G(y) is defined for every $y \in H = L^2(0, 1)$ by

$$G(y)(\xi) = g(\xi, y(\xi)).$$

In general, such functions are not Fréchet differentiable in *H*, but only Gâteaux differentiable, with the following expressions:

$$[DG(y).h](\xi) = \frac{\partial g}{\partial y}(\xi, y(\xi))h(\xi)$$
$$[D^2G(y).(h,k)](\xi) = \frac{\partial^2 g}{\partial y^2}(\xi, y(\xi))h(\xi)k(\xi)$$

We assume first that the linear coefficient *B* is given like in Example 2.3, with $H = L^2(0, 1)$. We remark that the finite dimensional subspaces H_N of *H* spanned by *N* eigenvectors of -B, with for k = 0, ..., N - 1

$$f_k(\xi) = \sqrt{2}\sin((k+1)\pi\xi),$$

are subspaces of the space of continuous functions $C([0, 1], \mathbb{R})$.

Then G_N is of class C^2 on the finite dimensional space H_N . The required estimates—holding uniformly with respect to N—on the second derivative are consequences of a Sobolev inequality: $D(-B)^{\eta}$ is continuously embedded into $L^{\infty}(0, 1)$ as soon as $\eta > 1/4$. For any h, k, ℓ , we have as a consequence of Holder inequality

$$\begin{aligned} \Big| &< D^2 G_N(y).(h,k), \ell > \Big| \le C \|hk\ell\|_{L^1(0,1)} \\ &\le C \sup \left(\|h\|_{L^2(0,1)} \|k\|_{L^2(0,1)} \|\ell\|_{L^\infty(0,1)}, \|h\|_{L^2(0,1)} \|k\|_{L^\infty(0,1)} \|\ell\|_{L^2(0,1)} \right) \\ &\le C \sup \left(|h|_H |k|_H |\ell|_{(-B)^\eta}, |h|_H |k|_{(-B)^\eta} |\ell|_H \right). \end{aligned}$$

See also [26] (Chapter 4).

2.3 The Cylindrical Wiener Process and Stochastic Integration in H

In this section, we recall the definition of the cylindrical Wiener process and of stochastic integral on a separable Hilbert space H with norm $|.|_H$. For more details, see [3].

We first fix a filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t\geq 0}, \mathbb{P})$. A cylindrical Wiener process on *H* is defined with two elements:

- a complete orthonormal system of H, denoted by (q_i)_{i∈I}, where I is a subset of N;
- a family (β_i)_{i∈I} of independent real Wiener processes with respect to the filtration ((F_t)_{t≥0});

then W is defined by

$$W(t) = \sum_{i \in I} \beta_i(t) q_i.$$
(4)

When *I* is a finite set, we recover the usual definition of Wiener processes in the finite dimensional space $\mathbb{R}^{|I|}$. However the subject here is the study of some Stochastic Partial Differential Equations, so that in the sequel the underlying Hilbert space *H* is infinite dimensional; for instance when $H = L^2(0, 1)$, an example of complete orthonormal system is $(q_k) = (\sqrt{2} \sin(k\pi .))_{k>1}$ —see Example 2.3.

A fundamental remark is that the series in Eq. 4 does not converge in H; but if a linear operator $\Psi : H \to K$ is Hilbert-Schmidt, then $\Psi W(t)$ converges in $L^2(\Omega, H)$ for any $t \ge 0$.

We recall that a bounded linear operator $\Psi : H \to K$ is said to be Hilbert-Schmidt when

$$|\Psi|^2_{\mathcal{L}_2(H,K)} := \sum_{k=0}^{+\infty} |\Psi(q_k)|^2_K < +\infty,$$

where the definition is independent of the choice of the orthonormal basis (q_k) of H. The space of Hilbert-Schmidt operators from H to K is denoted $\mathcal{L}_2(H, K)$; endowed with the norm $|.|_{\mathcal{L}_2(H,K)}$ it is an Hilbert space.

The stochastic integral $\int_0^t \Psi(s) dW(s)$ is defined in *K* for predictible processes Ψ with values in $\mathcal{L}_2(H, K)$ such that $\int_0^t |\Psi(s)|^2_{\mathcal{L}_2(H,K)} ds < +\infty$ a.s; moreover when $\Psi \in L^2(\Omega \times [0, t]; \mathcal{L}_2(H, K))$, the following two properties hold:

$$\mathbb{E} |\int_0^t \Psi(s) dW(s)|_K^2 = \mathbb{E} \int_0^t |\Psi(s)|_{\mathcal{L}_2(H,K)}^2 ds \text{ (Itô isometry),}$$
$$\mathbb{E} \int_0^t \Psi(s) dW(s) = 0.$$

A generalization of Itô formula also holds—see [3].

For instance, if $v = \sum_{k \in \mathbb{N}} v_k q_k \in H$, we can define

$$\langle W(t), v \rangle = \int_0^t \langle v, dW(s) \rangle = \sum_{k \in \mathbb{N}} \beta_k(t) v_k;$$

we then have the following space-time white noise property

$$\mathbb{E} < W(t), v_1 > < W(s), v_2 > = t \land s < v_1, v_2 > .$$

Therefore to be able to integrate a process with respect to *W* requires some strong properties on the integrand; in our SPDE setting, the Hilbert-Schmidt properties follow from the assumptions made on the linear coefficients of the equations.

Thanks to Assumption 2.2, it is easy to show that the following stochastic integral is well-defined in *H*, for any $t \ge 0$:

$$W^{B}(t) = \int_{0}^{t} e^{(t-s)B} dW(s).$$
 (5)

It is called a stochastic convolution, and it is the unique mild solution of

$$dZ(t) = BZ(t)dt + dW(t)$$
 with $Z(0) = 0$.

Under the second condition of Assumption 2.2, there exists $\delta > 0$ such that for any t > 0 we have $\int_0^t \frac{1}{s^\delta} |e^{sB}|^2_{\mathcal{L}_2(H)} ds < +\infty$; it can then be proved that W^B has continuous trajectories—via the *factorization method*, see [3]—and that for any $1 \le p < +\infty$

$$\mathbb{E}\sup_{t\geq 0}|W^B(t)|^p_H < +\infty.$$
(6)

We can now define solutions to Eq. 2, thanks to the assumptions made on the coefficients: the following result is classical—see [3]:

Proposition 2.8 For every T > 0, $y \in H$, the Eq. 2 admits a unique mild solution $Y \in L^2(\Omega, C([0, T], H))$:

$$Y(t) = e^{tB}y + \int_0^t e^{(t-s)B}G(Y(s))ds + \int_0^t e^{(t-s)B}dW(s).$$
 (7)

3 Definition of the Numerical Scheme

We now define the numerical approximation of Y: denoting by τ the time step, we have

$$Y_{k+1}(\tau, y) = Y_k(\tau, y) + \tau B Y_{k+1}(\tau, y) + \tau G(Y_k(\tau, y)) + \sqrt{\tau} \chi_{k+1}$$
$$Y_0(\tau, y) = y,$$

where $\chi_{k+1} = \frac{1}{\sqrt{\tau}} (W((k+1)\tau) - W(k\tau)).$

To simplify the equations, we omit the dependence of Y_k on the time-step τ and on the initial condition y.

This expression does not make sense in *H*. Defining $R_{\tau} = (I - \tau B)^{-1}$, this last equation can be replaced by

$$Y_{k+1} = R_{\tau} Y_k + \tau R_{\tau} G(Y_k) + \sqrt{\tau} R_{\tau} \chi_{k+1}, \qquad (8)$$

which is valid, since R_{τ} is a Hilbert-Schmidt operator on *H*.

Remark 3.1 Later, we often use the following expression for Y_k :

$$Y_k = R_{\tau}^k y + \tau \sum_{l=0}^{k-1} R_{\tau}^{k-l} G(Y_l) + \sqrt{\tau} \sum_{l=0}^{k-1} R_{\tau}^{k-l} \chi_{l+1}.$$
(9)

The following expression is also useful:

$$\sqrt{\tau} \sum_{l=0}^{k-1} R_{\tau}^{k-l} \chi_{l+1} = \int_0^{t_k} R_{\tau}^{k-l_s} dW(s), \tag{10}$$

where $l_s = \lfloor \frac{s}{\tau} \rfloor$ —with the notation $\lfloor . \rfloor$ for the integer part.

We need the following technical estimates:

Lemma 3.2 For any $0 \le \kappa \le 1$ and $j \ge 1$,

$$|(-B)^{1-\kappa}R^{j}_{\tau}|_{\mathcal{L}(H)} \leq c \frac{1}{(j\tau)^{1-\kappa}} \frac{1}{(1+\mu_{0}\tau)^{j\kappa}}$$

Moreover, for any $\beta \ge 1$ *there exists a constant* c_{β} *such that if* $j \in \mathbb{N}$ *with* $j \ge \beta$ *, then*

$$|(-B)^{\beta}R^{j}_{\tau}|_{\mathcal{L}(H)} \leq c_{\beta}\frac{1}{(j\tau)^{\beta}}.$$

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Proof For any $z \in H$,

$$\begin{split} |(-B)^{1-\kappa} R_{\tau}^{j} z|_{H}^{2} &= \sum_{i=0}^{+\infty} \mu_{i}^{2(1-\kappa)} \frac{1}{(1+\mu_{i}\tau)^{2j}} |z_{i}|^{2} \\ &= \frac{1}{(j\tau)^{2(1-\kappa)}} \sum_{i=0}^{+\infty} |z_{i}|^{2} \mu_{i}^{2(1-\kappa)} (j\tau)^{2(1-\kappa)} \frac{1}{(1+\mu_{i}\tau)^{2j(1-\kappa)}} \frac{1}{(1+\mu_{i}\tau)^{2j\kappa}} \\ &\leq \frac{1}{(j\tau)^{2(1-\kappa)}} \sum_{i=0}^{+\infty} \left(\frac{\mu_{i}j\tau}{1+\mu_{i}j\tau} \right)^{2(1-\kappa)} \frac{1}{(1+\mu_{0}\tau)^{2j\kappa}} |z_{i}|^{2} \\ &\leq c |z|_{H}^{2} \frac{1}{(j\tau)^{2(1-\kappa)}} \frac{1}{(1+\mu_{0}\tau)^{2j\kappa}}. \end{split}$$

The proof of the second inequality is similar: if $j \ge \beta$

$$\begin{split} |(-B)^{\beta} R_{\tau}^{j} z|_{H}^{2} &= \sum_{i=0}^{+\infty} \mu_{i}^{2\beta} \frac{1}{(1+\mu_{i}\tau)^{2j}} |z_{i}|^{2} = \sum_{i=0}^{+\infty} \mu_{i}^{2\beta} \frac{1}{(1+\mu_{i}\tau)^{2\beta\frac{j}{\beta}}} |z_{i}|^{2} \\ &\leq \sum_{i=0}^{+\infty} \mu_{i}^{2\beta} \frac{1}{(1+\mu_{i}\tau\frac{j}{\beta})^{2\beta}} |z_{i}|^{2} \leq \frac{\beta^{\beta}}{(j\tau)^{\beta}} |z|_{H}^{2}. \end{split}$$

4 Preliminary Results

We warn the reader that constants may vary from line to line during the proofs, and that in order to use lighter notations we usually forget to mention dependence on the parameters. We use the generic notation C for such constants.

We fix the time step τ , as well as $m \in \mathbb{N}$; we then introduce the notation $T = m\tau$. We also define $t_k = k\tau$. $\kappa > 0$ is a parameter, which is be supposed to be small enough. We also control τ : for some $\tau_0 > 0$, $\tau \leq \tau_0$.

4.1 Galerkin Approximation

The first step of the proof is to consider finite dimensional approximations of the *H*-valued processes $(Y(t))_{t \in \mathbb{R}^+}$ and $(Y_k)_{k \in \mathbb{N}}$: if we fix $N \ge 1$, we define $(Y^{(N)}(t))_{t \in \mathbb{R}^+}$ and $(Y_k^{(N)})_{k\in\mathbb{N}}$ by the equations

$$dY^{(N)}(t) = BY^{(N)}(t)dt + G_N(Y^{(N)}(t))dt + dW^{(N)}(t)$$

and

$$Y_{k+1}^{(N)} = Y_k^{(N)} + \tau B Y_{k+1}^{(N)} + \tau G_N(Y_k^{(N)}) + \sqrt{\tau} P_N \chi_{k+1}$$

with the initial conditions $Y_{t=0}^{(N)} = Y_{k=0}^{(N)} = P_N y$. The projection P_N and the nonlinear coefficient G_N have been defined above. $W^{(N)} = P_N W$ is a N-dimensional Wiener process on the subspace H_N . We remark that the above equations are well-defined on H_N —which is a stable subspace of B.

It is not difficult to prove the following convergence results: for any fixed $t \in \mathbb{R}^+$ and $k \in \mathbb{N}$, when $N \to +\infty$ we have

$$\mathbb{E}|Y(t) - Y^{(N)}(t)|^2 \to 0 \text{ and } \mathbb{E}|Y_k - Y_k^{(N)}|^2 \to 0.$$

For this, we use the following inequality—where *I* denotes the identity map on *H*: for any $0 < \beta < 1$

$$|(-B)^{-\beta}(I-P_N)|_{\mathcal{L}(H)} \le \frac{1}{\mu_{N+1}^{\beta}}.$$

Given 0 < r < 1/2, we then obtain for any $t \ge 0$, $k \in \mathbb{N}$ and $\tau > 0$

$$\begin{split} \mathbb{E} |\int_{0}^{t} e^{(t-s)B} (I-P_{N}) dW(s)|^{2} \\ &= \int_{0}^{t} \operatorname{Tr} \left((I-P_{N}) e^{(t-s)B} (-B)^{-1/2-r} (-B)^{1-r} (-B)^{-1/2+2r} e^{(t-s)B} (I-P_{N}) \right) ds \\ &\leq \operatorname{Tr} ((-B)^{-1/2-r}) |(-B)^{-1/2+2r} (I-P_{N})|_{\mathcal{L}(H)} \int_{0}^{t} \frac{C_{r}}{(t-s)^{1-r}} ds, \end{split}$$

and the last expression goes to 0 when $N \to +\infty$. To conclude, it remains to use the Gronwall lemma. The proof for the discrete-time process is similar.

We need test functions Φ_N adapted to the finite-dimensional approximation: for any $N \ge 1$, by restriction we define $\Phi_N(y) = \Phi(y)$ for any $y \in H_N$; we obtain the following decomposition

(1)

$$\mathbb{E}\Phi(Y(m\tau)) - \mathbb{E}\Phi(Y_m) = \mathbb{E}\Phi(Y(m\tau)) - \mathbb{E}\Phi(Y^{(N)}(m\tau)) + \mathbb{E}\Phi_N(Y^{(N)}(m\tau)) - \mathbb{E}\Phi_N(Y_m^{(N)}) + \mathbb{E}\Phi(Y_m^{(N)}) - \mathbb{E}\Phi(Y_m);$$

the first and the third terms converge to 0 when $N \to +\infty$. In the sequel, we prove an estimate of the second term, which is uniform with respect to dimension N; letting $N \to +\infty$ then yields an estimate on the left hand side.

Hence we work with the finite dimensional approximation, but we omit the parameter N. The constants appearing below are independent of N.

In Section 4.2, we prove some estimates on Y(t) and Y_m , and in Section 4.3 we focus on the asymptotic behaviour of the processes.

4.2 Some Useful Estimates

Bounds on moments of Y_t and Y_k can be proved, uniformly with respect to time.

Lemma 4.1 For any $p \ge 1$, there exists a constant $C_p > 0$ such that for every $t \ge 0$ and $y \in H$

$$\mathbb{E}|Y(t, y)|^p \le C_p(1+|y|^p).$$

Proof If we define $Z(t) = Y(t) - W^B(t)$, we have Z(0) = Y(0) = y, and

$$\frac{dZ(t)}{dt} = BZ(t) + G(Y(t)),$$

and by Proposition 2.6

$$\begin{aligned} \frac{1}{2} \frac{d|Z(t)|^2}{dt} &= \langle BZ(t) + G(Y(t)), Z(t) \rangle \\ &= \langle BZ(t) + G(Z(t)), Z(t) \rangle + \langle G(Y(t)) - G(Z(t)), Z(t) \rangle \\ &\leq -c|Z(t)|^2 + C + \|G\|_{\infty} |Z(t)| \\ &\leq -c'|Z(t)|^2 + C', \end{aligned}$$

for some new constants c', C'.

Then almost surely we have for any $t \ge 0$

$$|Z(t)| \le C(1+|y|).$$

Thanks to Eq. 6, the conclusion easily follows.

Lemma 4.2 For any $p \ge 1$, $\tau_0 > 0$, there exists a constant C > 0 such that for every $0 < \tau \le \tau_0$, $k \in \mathbb{N}$ and $y \in H$

$$\mathbb{E}|Y_k|^p \le C(1+|y|^p).$$

Proof As in the proof of Lemma 4.1 above, we introduce $Z_m = Y_m - w_m$, where the process (w_m) is the numerical approximation of W^B with the numerical scheme (8)—with G = 0; it is defined by

$$w_{m+1} = R_\tau w_m + \sqrt{\tau} R_\tau \chi_{m+1}.$$

Using Theorem 3.2 of [25], giving the strong order 1/4 for the numerical scheme when the initial condition is 0, with no nonlinear coefficient, with a constant diffusion term and under the assumptions made here—we obtain the following estimate: for any $p \ge 1$, $\tau_0 > 0$ and 0 < r < 1/2 there exists C > 0 such that for any $0 < \tau \le \tau_0$ and $m \ge 0$

$$\mathbb{E}|w_m - W^B(m\tau)|^{2p} \le C\tau^{(1/2-r)p}.$$
(11)

Thanks to Eqs. 6 and 11, we get that for any $\tau_0 > 0$, there exists C > 0 such that for $0 < \tau \le \tau_0$ and $m \ge 0$

$$4\mathbb{E}|w_m|^2 \le C. \tag{12}$$

Now Z_m defined above satisfies $Z_0 = Y_0 = y$ and

$$Z_{m+1} = R_{\tau} Z_m + \tau R_{\tau} G(Y_m);$$

since $|R_{\tau}|_{\mathcal{L}(H)} \leq \frac{1}{1+\mu_0\tau}$, we obtain the almost sure estimates

$$|Z_{m+1}| \le \frac{1}{1 + \mu_0 \tau} |Z_m| + C\tau$$

and

$$|Z_m| \le C(1+|y|).$$

Thanks to Eq. 12, we therefore obtain the result.

We now introduce the following process: for $0 \le k \le m - 1$ and $t_k \le t \le t_{k+1}$

$$\tilde{Y}(t) = Y_k + \int_{t_k}^t [B_\tau Y_k + R_\tau G(Y_k)] ds + \int_{t_k}^t R_\tau dW(s),$$
(13)

where $B_{\tau} = BR_{\tau}$. The process \tilde{Y} is a natural interpolation of the numerical solution (Y_k) defined by Eq. 8: $\tilde{Y}(t_k) = Y_k$.

Thanks to Lemma 4.2, we get

Lemma 4.3 For any $p \ge 1$, $\tau_0 > 0$, there exists C > 0 such that for any $0 < \tau \le \tau_0$, $t \ge 0$ and $y \in H$

$$\mathbb{E}|\tilde{Y}(t)|^p \le C(1+|y|^p).$$

In the next Lemma, we give a control on Malliavin derivatives of Y_k used in the proof. For an introduction to Malliavin calculus, see [24, 27]. We take the notation \mathcal{D} for the Malliavin derivative, and we follow the presentation given in [6], where the following useful integration by parts formula is given—see Lemma 2.1 therein:

Lemma 4.4 For any $F \in \mathbb{D}^{1,2}(H)$, $u \in C^2(H)$ with bounded derivatives and $\Psi \in L^2(\Omega \times [0, T], \mathcal{L}_2(H))$ an adapted process,

$$\mathbb{E}\left[Du(F) \int_{0}^{T} \Psi(s) dW(s)\right] = \mathbb{E}\left[\int_{0}^{T} Tr(\Psi(s)^{*}D^{2}u(F)\mathcal{D}_{s}F) ds\right],$$
(14)

where $\mathcal{D}_s F : h \in H \mapsto \mathcal{D}_s^h F \in H$ stands for the Malliavin derivative of F, and $\mathbb{D}^{1,2}(H)$ is the set of H-valued random variables $F = \sum_{i \in \mathbb{N}} F_i f_i$, with $F_i \in \mathbb{D}^{1,2}$ the domain of the Malliavin derivative for \mathbb{R} -valued random variables for any i, and $\sum_{i \in \mathbb{N}} \int_0^T \mathbb{E} |\mathcal{D}_s F_i|^2 ds < +\infty$.

Without any stronger dissipativity assumption, we are not able to give a uniform control with respect to time of the Malliavin derivative of \tilde{Y} . In the proof below, we take care of this problem by using these derivatives only at times $t_k = k\tau$ and s such that $t_{k-l_s} \leq 1$.

Lemma 4.5 For any $0 \le \beta < 1$ and $\tau_0 > 0$, there exists a constant C > 0 such that for every $h \in H$, $k \ge 1$, $0 < \tau \le \tau_0$ and $s \in [0, t_k]$

$$|\mathcal{D}_{s}^{h}Y_{k}|_{\beta} \leq C(1+L_{G}\tau)^{k-l_{s}}\left(1+\frac{1}{(1+\mu_{0}\tau)^{(1-\beta)(k-l_{s})}t_{k-l_{s}}^{\beta}}\right)|h|$$

For any $k \ge 1, h \in H$ and $s \in [0, t_k]$, using the chain rule for Malliavin calculus and expressions (9) and (10), we have

$$\mathcal{D}_s^h Y_k = R_\tau^{k-l_s} h + \tau \sum_{l=l_s+1}^{k-1} R_\tau^{k-l} DG(Y_l) \cdot \mathcal{D}_s^h Y_l$$

Indeed, recall that l_s denotes the integer part of $\frac{s}{\tau}$, so that when $l \leq l_s$ we have $\mathcal{D}_s^h Y_l = 0$.

As a consequence, for $k \ge l_s + 1$

$$|\mathcal{D}_s^h Y_k| \le (1 + L_G \tau)^{k - l_s} |h|$$

Now using Lemma 3.2, we have

$$\begin{aligned} |(-B)^{\beta} \mathcal{D}_{s}^{h} Y_{k}| &\leq \frac{1}{(1+\mu_{0}\tau)^{(1-\beta)(k-l_{s})} t_{k-l_{s}}^{\beta}} |h| \\ &+ L_{G} \tau \sum_{l=l_{s}+1}^{k-1} \frac{(1+L_{G}\tau)^{l-l_{s}}}{(1+\mu_{0}\tau)^{(1-\beta)(k-l)} t_{k-l}^{\beta}} |h|. \end{aligned}$$

To conclude, we see that

$$\tau \sum_{l=l_{s}+1}^{k-1} \frac{1}{(1+\mu_{0}\tau)^{(1-\beta)(k-l)} t_{k-l}^{\beta}} \le C \int_{0}^{+\infty} t^{-\beta} \frac{1}{(1+\mu_{0}\tau)^{(1-\beta)\frac{L}{\tau}}} dt \le C < +\infty,$$

when $0 < \tau \leq \tau_0$.

4.3 Asymptotic Behaviour of the Processes

The results of this section are obtained for the initial *H*-valued processes, and for their finite dimensional approximations.

First, we focus on the existence of invariant measures for the continuous and discrete time processes. We use the well-known Krylov-Bogoliubov criterion—see [4]. Tightness comes from two facts: $D(-B)^{\gamma}$ is compactly embedded in H when $\gamma > 0$, and when $\gamma < 1/4$ we can control moments:

Lemma 4.6 For any $0 < \gamma < 1/4$, $\tau > 0$ and any $y \in H$, there exists $C(\gamma, \tau, y)$, $C(\gamma, y) > 0$ such that for any $m \ge 1$ and $t \ge 1$

$$\mathbb{E}|Y_m(\tau, y)|_{\gamma}^2 \le C(\gamma, \tau, y) \quad and \quad \mathbb{E}|Y(t, y)|_{\gamma}^2 \le C(\gamma, y)$$

Uniqueness of the invariant probability measure for the continuous time process $(Y(t))_{t \in \mathbb{R}^+}$ can be deduced from the well-known Doob Theorem—see [4]. Indeed, since in Eq. 2 noise is additive and non-degenerate, the Strong Feller property—see also Lemma 5.6 below—and irreducibility can be easily proved. In the proof of the main Theorem 1.1, we also need speed of convergence, and thanks to a coupling argument we get the following exponential convergence result—for a proof see Theorem 2.4 and Section 6.1 in [7]:

Proposition 4.7 There exist c > 0, C > 0 such that for any bounded test function ϕ , any $t \ge 0$ and any $y_1, y_2 \in H$

$$|\mathbb{E}\phi(Y(t, y_1)) - \mathbb{E}\phi(Y(t, y_2))| \le C \|\phi\|_{\infty} (1 + |y_1|^2 + |y_2|^2) e^{-ct}.$$
(15)

The idea of coupling relies on the following formula: if v_1 and v_2 are two probability measures on a state space *S*, their total variation distance satisfies

$$d_{TV}(v_1, v_2) = \inf \{\mathbb{P}(X_1 \neq X_2)\}$$

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which is an infimum over random variables (X_1, X_2) defined on a same probability space, and such that $X_1 \sim \nu_1$ and $X_2 \sim \nu_2$.

Roughly speaking, the principle is to define a coupling $(Z_1(t, y_1, y_2), Z_2(t, y_1, y_2))_{t\geq 0}$ for the processes $(Y(t, y_1)_{t\geq 0}$ and $Y((t, y_2))_{t\geq 0}$ such that the coupling time T of Z_1 and Z_2 —i.e. the first time the processes are equal—has an exponentially decreasing tail.

This technique was first used in the study of the asymptotic behaviour of Markov chains—see [1, 9, 16, 20]—and was later adapted for SDEs and more recently for SPDEs—see for instance [15, 18, 23].

In fact, uniqueness of an invariant probability measure $\overline{\mu}$ is an easy consequence of this Proposition, and moreover we get for any $y \in H$ and any $t \ge 0$

$$\left| \mathbb{E}\phi(Y(t,y)) - \int_{H} \phi d\overline{\mu} \right| \le C \|\phi\|_{\infty} (1+|y|^{2}) e^{-ct}.$$
 (16)

In general, we do not know whether uniqueness also holds for the numerical approximation $(Y_k(\tau, .))_{k \in \mathbb{N}}$.

Remark 4.8 A sufficient condition for the uniqueness of the invariant probability measure of the discrete time process $(Y_k)_{k \in \mathbb{N}}$ is the strict dissipativity assumption

$$L_G < \mu_0,$$

where we recall that L_G denotes the Lipschitz constant of G.

Then trajectories of the processes $(Y_t)_{t \in \mathbb{R}^+}$ and $(Y_k)_{k \in \mathbb{N}}$ issued from different initial conditions y_1 and y_2 and driven by the same noise process are exponentially close when time increases: for any $\tau_0 > 0$, there exists c > 0 such that for any $0 < \tau \le \tau_0$, $k \ge 0$ and $t \ge 0$ we have almost surely

$$\begin{aligned} |Y(t, y_1) - Y(t, y_2)| &\leq e^{-(\mu_0 - L_G)t} |y_1 - y_2| \\ |Y_k(\tau, y_1) - Y_k(\tau, y_2)| &\leq e^{-ck\tau} |y_1 - y_2|. \end{aligned}$$

Proof of uniqueness is then straightforward—and we do not need Proposition 4.7 above.

5 Presentation of the Proof of the Weak Approximation Result

The proof of Theorem 1.1 is very technical, so for pedagogy we first intro*D*uce the decomposition of the error, and identify the term which we control later in Section 6. Some crucial estimates on the derivatives of the semi-group with respect to the initial conditions—regularization, long-time behaviour—are proved below in Section 5.2.

5.1 Strategy

We define

$$u(t, y) = \mathbb{E}[\phi(Y(t, y))], \tag{17}$$

which is solution of a finite dimensional Kolmogorov equation associated with the finite dimensional approximation of Eq. 2:

$$\frac{du}{dt}(t, y) = Lu(t, y) = \frac{1}{2} \text{Tr} \left(D^2 u(t, y) \right) + \langle By + G(y), Du(t, y) \rangle$$

As explained in the introduction, this is one of the essential tools in the proof of the weak approximation result.

The weak error at time $T = m\tau$ can be decomposed with a telescoping sum where to simplify the dependence of the numerical approximation in τ and y is not written:

$$\mathbb{E}[\phi(Y(T, y))] - \mathbb{E}[\phi(Y_m)] = u(T, y) - \mathbb{E}[u(0, Y_m)]$$

$$= \sum_{k=0}^{m-1} (\mathbb{E}[u(T - t_k, Y_k)] - \mathbb{E}[u(T - t_{k+1}, Y_{k+1})])$$

$$= u(T, y) - \mathbb{E}[u(T - \tau, Y_1(\tau, y))] + \sum_{k=1}^{m-1} (a_k + b_k + c_k),$$
(18)

where for $1 \le k \le m - 1$

$$a_{k} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} \langle B\tilde{Y}(t) - B_{\tau}Y_{k}, Du(T-t, \tilde{Y}(t)) \rangle dt,$$

$$b_{k} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} \langle G(\tilde{Y}(t)) - R_{\tau}G(Y_{k}), Du(T-t, \tilde{Y}(t)) \rangle dt,$$

$$c_{k} = \frac{1}{2}\mathbb{E} \int_{t_{k}}^{t_{k+1}} \operatorname{Tr}((I - R_{\tau}R_{\tau}^{*})D^{2}u(T-t, \tilde{Y}(t)))dt.$$
(19)

This follows from the use of the Kolmogorov equation and of the Itô formula. We recall that \tilde{Y} is defined in Eq. 13.

5.2 Bounds on the Derivatives of the Transition Semi-group

By Eq. 17, $u(t, y) = \mathbb{E}[\phi(Y(t, y))]$; since ϕ is of class C^2 , bounded and with bounded derivatives, we are able to prove that with respect to *y* the function *u* is twice differentiable, and that the derivatives can be calculated in the following way:

• For any $h \in H$, we have

$$Du(t, y).h = \mathbb{E}[D\phi(Y(t, y)).\eta^{h, y}(t)],$$
(20)

where $\eta^{h,y}$ is the solution of

$$\frac{d\eta^{h,y}(t)}{dt} = B\eta^{h,y}(t) + DG(Y(t, y)).\eta^{h,y}(t),$$
$$\eta^{h,y}(0) = h.$$

• For any $h, k \in H$, we have

$$D^{2}u(t, y).(h, k) = \mathbb{E}[D^{2}\phi(Y(t, y)).(\eta^{h, y}(t), \eta^{k, y}(t)) + D\phi(Y(t, y)).\zeta^{h, k, y}(t)],$$
(21)

where $\zeta^{h,k,y}$ is the solution of

$$\frac{d\zeta^{h,k,y}(t)}{dt} = B\zeta^{h,k,y}(t) + DG(Y(t,y)).\zeta^{h,k,y}(t) + D^2G(Y(t,y)).(\eta^{h,y}(t),\eta^{k,y}(t)),$$

$$\zeta^{h,k,y}(0) = 0.$$

In [6], the key point for obtaining the weak order 1/2 is to control the derivatives $|Du(t, y)|_{\beta}$ and $|(-B)^{\beta}D^{2}u(t, y)(-B)^{\gamma}|_{\mathcal{L}(H)}$, with $\beta < 1/2$ and $\gamma < 1/2$ —with the identifications of Remark 2.1. Moreover, to obtain a long-time weak estimate we need to prove some exponential decreasing of such quantities when time *t* goes to infinity. The two Propositions below are the essential results we thus need.

Proposition 5.1 *There exists a constant* $\tilde{\mu} > 0$ *such that for any* $0 \le \beta < 1$ *, for any* t > 0 *and* $y \in H$

$$|Du(t, y)|_{\beta} \le C_{\beta} \left(1 + \frac{1}{t^{\beta}} \right) e^{-\tilde{\mu}t} (1 + |y|^2).$$
(22)

Proposition 5.2 There exists a constant $\tilde{\mu} > 0$ such that for any $0 \le \beta, \gamma < 1/2$, for any t > 0 and $y \in H$

$$|(-B)^{\beta} D^{2} u(t, y)(-B)^{\gamma}|_{\mathcal{L}(H)} \le C_{\beta, \gamma} \left(1 + \frac{1}{t^{\eta}} + \frac{1}{t^{\beta+\gamma}}\right) e^{-\tilde{\mu}t} (1 + |y|^{2}).$$
(23)

Remark 5.3 According to the identifications made in Remark 2.1, Proposition 5.1 is a bound on the norm of $Du(t, y) \in H$, with respect to the norm $|.|_{\beta}$.

Proposition 5.2 is a bound on the linear operator $(-B)^{\beta}D^{2}u(t, y)(-B)^{\gamma}$, defined such that for any $h, k \in H$

$$< (-B)^{\beta} D^{2} u(t, y) (-B)^{\gamma} h, k > = < D^{2} u(t, y) \left((-B)^{\beta} h \right), (-B)^{\gamma} k >$$
$$= D^{2} u(t, y) \left((-B)^{\beta} h, (-B)^{\gamma} h \right),$$

where in the last equality we just rewrite the identification of the second-order derivative with a linear operator.

The singularity $t^{-\eta}$ in Eq. 23 is a consequence of the regularity properties satisfied by G. Since in general during the proof of Theorem 1.1, we need $\beta + \gamma$ to be close to 1, and therefore greater than η , only the second singularity $t^{-\beta-\gamma}$ plays a role.

The proofs require several steps. First in Lemma 5.4 below we prove estimates for a finite horizon and general $0 \le \beta$, $\gamma < 1/2$; then in Lemma 5.6 we study the long-time behaviour in the particular case $\beta = \gamma = 0$; we finally conclude with the proofs of Propositions 5.1 and 5.2.

First, we prove estimates of these quantities for $0 < t \le 1$ —see Lemmas 4.4 and 4.5 in [6], with a difference coming from the assumptions made on the nonlinear coefficient *G*:

Lemma 5.4 For any $0 \le \beta < 1/2$, $0 \le \gamma < 1/2$, there exist constants C_{β} and $C_{\beta,\gamma}$ such that for any $y \in H$ and any $0 < t \le 1$

$$\begin{aligned} |Du(t,y)|_{\beta} &\leq \frac{C_{\beta}}{t^{\beta}} \\ |(-B)^{\beta} D^{2} u(t,y)(-B)^{\gamma}|_{\mathcal{L}(H)} &\leq C_{\beta,\gamma} \left(\frac{1}{t^{\eta}} + \frac{1}{t^{\beta+\gamma}}\right). \end{aligned}$$

Remark 5.5 If we take another time interval $]0, T_{max}]$ instead of]0, 1], the constants C_{β} and $C_{\beta,\gamma}$ are *a priori* exponentially increasing in T_{max} .

The estimate on *Du* is valid for $\beta < 1$, and not only for $\beta < 1/2$.

Proof Owing to Eqs. 20 and 21, we only need to prove the following almost sure estimates, for some constants C_{β} and $C_{\beta,\gamma}$ —which may vary from line to line below: for any $0 < t \le 1$

$$\begin{aligned} |\eta^{h,y}(t)| &\leq \frac{C_{\beta}}{t^{\beta}} |h|_{-\beta} \\ |\zeta^{h,k,y}(t)| &\leq \frac{C_{\beta,\gamma}}{t^{\eta}} |h|_{-\beta} |k|_{-\gamma}, \end{aligned}$$
(24)

where the parameter η is defined in Assumption 2.5.

We use mild formulations, and the regularization properties of the semi-group given in Proposition 2.4:

$$\begin{aligned} |\eta^{h,y}(t)| &= \left| e^{tB}h + \int_0^t e^{(t-s)B} DG(Y(s,y)) \cdot \eta^{h,y}(s) ds \right| \\ &\leq \frac{C_\beta}{t^\beta} |h|_{-\beta} + C \int_0^t |\eta^{h,y}(s)| ds, \end{aligned}$$

and by the Gronwall Lemma we get the result.

For the second-order derivative, we moreover use the properties of G in Assumption 2.5 to get

$$\begin{aligned} |\zeta^{h,k,y}(t)| &= |\int_0^t e^{(t-s)B} DG(Y(s,y)).\zeta^{h,k,y}(s)ds \\ &+ \int_0^t e^{(t-s)B} D^2 G(Y(s,y)).(\eta^{h,y}(s),\eta^{k,y}(s))ds | \\ &\leq C \int_0^t |\zeta^{h,k,y}(s)|ds + \int_0^t \frac{C_{\beta,\gamma}}{(t-s)^{\eta}} |\eta^{h,y}(s)| |\eta^{k,y}(s)|ds \\ &\leq C \int_0^t |\zeta^{h,k,y}(s)|ds + C_{\beta,\gamma}|h|_{-\beta} |k|_{-\gamma} t^{1-\eta-\beta-\gamma} \int_0^1 \frac{1}{(1-s)^{\eta} s^{\beta+\gamma}} ds. \end{aligned}$$

To conclude, it remains to use the Gronwall Lemma, since for any $0 < t \le 1$ we get $t^{1-\eta-\beta-\gamma} \le t^{-\eta}$, thanks to the assumption $\beta + \gamma < 1$.

Thanks to the dissipativity property expressed in Proposition 2.6, we can prove the result in the case $\beta = \gamma = 0$. We notice that the proof would be straightforward under a strict dissipativity assumption—since then $\eta^{h,y}(t)$ and $\zeta^{h,k,y}(t)$ would decrease exponentially in *t*; in the general case $\eta^{h,y}(t)$ and $\zeta^{h,k,y}(t)$ are exponentially increasing in time so that we can not work directly. Here the result comes from the estimate (15) of Proposition 4.7.

Lemma 5.6 There exist constants C and c > 0 such that for any $t \ge 0$ and any $y \in H$

$$|Du(t, y)| \le Ce^{-ct}(1+|y|^2) \quad and \quad |D^2u(t, y)|_{\mathcal{L}(H)} \le Ce^{-ct}\left(1+\frac{1}{t^{\eta}}\right)(1+|y|^2). \tag{25}$$

Proof The Bismut-Elworthy-Li formula states that if $\Phi : H \to \mathbb{R}$ is a function of class C^2 with bounded derivatives and with at most quadratic growth—i.e. there exists $M(\Phi) > 0$ such that for any $y \in H$ we have $|\Phi(y)| \le M(\Phi)(1+|y|^2)$ —then we can calculate the first and the second order derivatives of $(t, y) \mapsto v(t, y) := \mathbb{E}\Phi(Y(t, y))$ with respect to y. First, we have for any $y \in H$ and $h \in H$

$$Dv(t, y).h = \frac{1}{t} \mathbb{E}\left[\int_0^t <\eta^{h, y}(s), dW(s) > \Phi(Y(t, y))\right]$$
$$= \frac{2}{t} \mathbb{E}\left[\int_0^{t/2} <\eta^{h, y}(s), dW(s) > v(t/2, Y(t/2, y))\right];$$
(26)

the second equality is a consequence of the identity $v(t, y) = \mathbb{E}v(t/2, Y(t/2, y))$ obtained with the Markov property, and of the first equality applied with the function v(t/2, .).

Using the second formula of Eq. 26, we obtain a formula for the second order derivative: for any $y \in H$ and $h, k \in H$,

$$D^{2}v(t, y).(h, k) = \frac{2}{t} \mathbb{E} \left[\int_{0}^{t/2} \langle \zeta^{h,k,y}(s), dW(s) \rangle v(t/2, Y(T/2, y)) \right] + \frac{2}{t} \mathbb{E} \left[\int_{0}^{t/2} \langle \eta^{h,y}(s), dW(s) \rangle Dv(t/2, Y(t/2)).\eta^{k,y}(t/2) \right].$$
(27)

We then see, using Lemmas 4.1 and 5.4—with $\beta = \gamma = 0$ —that there exists C > 0 such that for any $0 < t \le 1, y \in H, h, k \in H$

$$|Dv(t, y).h| \le \frac{C}{\sqrt{t}} M(\Phi)(1+|y|^2)|h|,$$

$$|D^2v(t, y).(h, k)| \le \frac{C}{t} M(\Phi)(1+|y|^2)|h||k|.$$
 (28)

Now when $t \ge 1$ the Markov property implies that $u(t, y) = \mathbb{E}u(t - 1, Y(1, y))$, and by Eq. 16 we have

$$|u(t-1, y) - \int_{H} \phi d\overline{\mu}| \le C e^{-c(t-1)} (1+|y|^{2}).$$

If we choose $\Phi_t(y) = u(t-1, y) - \int_H \phi d\overline{\mu}$, we have $u(t, y) = \mathbb{E}\Phi_t(Y(1, y)) + \int_H \phi d\overline{\mu}$, with $M(\Phi_t) \le Ce^{-c(t-1)}$. With Eq. 28 at time 1, we obtain for $t \ge 1$

$$|Du(t, y).h| \le Ce^{-c(t-1)}(1+|y|^2)|h|$$
$$|D^2u(t, y).(h, k)| \le Ce^{-c(t-1)}(1+|y|^2)|h||k|.$$

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Moreover by Lemma 5.4 we have a control when $0 \le t \le 1$, so that with a change of constants we get the result.

We can finally prove the Propositions 5.1 and 5.2. The key tool is the Markov property of the process *Y* which yields the following formula: for any $t \ge 1$

$$u(t, y) = \mathbb{E}[u(t-1, Y_1(y))].$$
(29)

To get the exponential decreasing, we use Lemma 5.6 at time t - 1 when $t \ge 1$, while $|h|_{-\beta}$ appears from $\eta_{h,v}(1)$, and with estimates coming from Lemma 5.4.

Proof of Propositions 5.1 and 5.2 Using Eq. 29 and Lemma 5.6, for any $t \ge 1$ we have

$$|Du(t, y).h| \le Ce^{-c(t-1)} \mathbb{E}[(1+|Y(1, y)|^2)|\eta^{h, y}(1)|] \le Ce^{-c(t-1)}(1+|y|^2)|h|_{-\beta},$$

where the last estimate comes from Lemmas 4.1 and 5.4. Combining this estimate with the result of Lemma 5.4, which gives an estimate for $t \le 1$, we obtain Eq. 22. For the second order derivatives, Lemma 5.4 gives an estimate for $t \le 1$, and for $t \ge 1$ we use Eq. 29 to see that

$$D^{2}u(t, y).(h, k) = \mathbb{E}[D^{2}[u(t-1, Y(1, y))].(h, k)]$$

= $\mathbb{E}D^{2}u(t-1, Y(1, y)).(\eta^{h, y}(1), \eta^{k, y}(1)) + \mathbb{E}Du(t-1, Y(1, y)).\zeta^{h, k, y}(1).$

Using Lemma 5.6, we get an exponential decreasing; thanks to Lemma 4.1 and to the estimates in the proof of Lemma 5.4 at time 1, we obtain

$$|D^{2}u(t, y).(h, k)| \le Ce^{-c(t-1)}(1+|y|^{2})|h|_{-\beta}|k|_{-\gamma}.$$

Then Eq. 23 easily follows.

5.3 Proof of Corollary 1.2

The first estimate is a simple consequence of Theorem 1.1, and of the exponential convergence to equilibrium of the continuous-time process—see Eq. 16. We then get

$$|\mathbb{E}[\phi(Y_m(\tau, y))] - \int_H \phi d\overline{\mu}| \le C(1+|y|^3) \left(\frac{1}{m^{1/2-\kappa}} + \tau^{1/2-\kappa}\right) + C(1+|y|^2)e^{-cm\tau}.$$

If μ^{τ} is an ergodic invariant probability measure of $(Y_m(\tau, .))_m$, then since ϕ is bounded for μ^{τ} -almost any $y \in H$ we have by the ergodic Theorem the following convergence when $M \to +\infty$:

$$\frac{1}{M}\sum_{m=1}^{M}\mathbb{E}[\phi(Y_m(\tau,y))] \to \int_{H}\phi(y)\mu^{\tau}(dy)$$

To conclude, it remains to choose a initial condition *y* in this non-empty set, and to use Cesaro Lemma on the right-hand side of the estimate.

We notice that if μ^{τ} is an invariant probability measure, not necessarily ergodic, having a finite moment of order 3, then it is enough to integrate the inequality above with respect to μ^{τ} .

6 Proof of the Estimates

We need to control the terms given in Eq. 19, according to the decomposition (18). We recall that constants C must be independent from the dimension N and the final time $T = m\tau$.

6.1 Estimate of $u(T, x) - \mathbb{E}[u(T - \tau, Y_1)]$

The Markov property gives

$$u(T, y) = \mathbb{E}[\phi(Y(T, y))] = \mathbb{E}[u(T - \tau, Y(\tau, y))].$$

If $0 < \kappa < 1/2$, using Lemma 4.2 and Proposition 5.1 we get

$$\begin{aligned} &|u(T, y) - \mathbb{E}[u(T - \tau, Y_1)]| \\ &\leq C(1 + (T - \tau)^{-1/2 + \kappa}) e^{-\tilde{\mu}(T - \tau)} (\mathbb{E}|Y(\tau, y) - Y_1|_{-1/2 + \kappa}^2)^{1/2} (1 + |y|^2). \end{aligned}$$

We can write

$$Y(\tau, y) - Y_1 = (e^{\tau B} - R_{\tau})y + \int_0^{\tau} e^{(\tau - s)B} G(Y(s, y)) ds - \tau R_{\tau} G(y) + \int_0^{\tau} e^{(\tau - s)B} dW(s) - \sqrt{\tau} R_{\tau} \chi_1.$$

We use the following properties to estimate the first line in this equality:

$$\begin{aligned} |(-B)^{-1/2+\kappa} (e^{\tau B} - R_{\tau})|_{\mathcal{L}(H)} &\leq c\tau^{1/2-\kappa}, \\ |e^{sB}|_{\mathcal{L}(H)} &\leq 1 \text{ for } s \geq 0, \\ &|R_{\tau}|_{\mathcal{L}(H)} \leq 1, \\ &|(-B)^{-1/2+\kappa}.| < c|.|; \end{aligned}$$

therefore the first line in the last expression is almost surely bounded by $C(\tau^{1/2-\kappa} + \tau)(1+|y|)$.

For the second line, we have

$$\mathbb{E} \left| (-B)^{-1/2+\kappa} \int_0^\tau e^{(\tau-s)B} dW(s) \right|^2 = \mathbb{E} \int_0^\tau |(-B)^{-1/2+\kappa} e^{(\tau-s)B}|^2_{\mathcal{L}_2(H)} ds$$

$$\leq \tau |(-B)^{-1/2+\kappa}|^2_{\mathcal{L}_2(H)}$$

$$\leq c\tau;$$

the last term is controlled in the same way: $\mathbb{E}|(-B)^{-1/2+\kappa}\sqrt{\tau}R_{\tau}\chi_1|^2 \le c\tau$. Therefore we have

$$|u(T, x) - \mathbb{E}[u(T - \tau, Y_1)]| \le C(1 + |y|^3)(1 + (T - \tau)^{-1/2 + \kappa})e^{-\tilde{\mu}(T - \tau)\tau^{1/2 - \kappa}}.$$
 (30)

We thus understand that to obtain weak order 1/2 requires to be careful in the estimate. Here we used Lemma 4.2 instead of Lemma 5.6; otherwise looking at $\mathbb{E}|Y(\tau, y) - Y_1|^2$ would have not been sufficient. The control of the other terms must be done in the same spirit.

6.2 Estimate of a_k

We have

$$\begin{aligned} a_k &= \mathbb{E} \int_{t_k}^{t_{k+1}} < B\tilde{Y}(t) - B_\tau Y_k, Du(T - t, \tilde{Y}(t)) > dt \\ &= \mathbb{E} \int_{t_k}^{t_{k+1}} < (B - B_\tau) Y_k, Du(T - t, \tilde{Y}(t)) > dt \\ &+ \mathbb{E} \int_{t_k}^{t_{k+1}} < B(\tilde{Y}(t) - Y_k), Du(T - t, \tilde{Y}(t)) > dt \\ &:= a_k^1 + a_k^2. \end{aligned}$$

6.2.1 Estimate of a_k^1

We use the equality $B_{\tau} - B = \tau R_{\tau} B^2$. We also decompose a_k^1 using expression (9):

$$\begin{aligned} a_k^{1,1} &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < R_\tau B^2 R_\tau^k y, Du(T-t, \tilde{Y}(t)) > dt \\ a_k^{1,2} &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < R_\tau B^2 \tau \sum_{l=0}^{k-1} R_\tau^{k-l} G(Y_l), Du(T-t, \tilde{Y}(t)) > dt \\ a_k^{1,3} &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < R_\tau B^2 \sqrt{\tau} \sum_{l=0}^{k-1} R_\tau^{k-l} \chi_{l+1}, Du(T-t, \tilde{Y}(t)) > dt; \end{aligned}$$

then $a_k^1 = a_k^{1,1} + a_k^{1,2} + a_k^{1,3}$.

 Estimate of a_k^{1,1}
 The idea is to "share" B² between different factors—thanks to regularization properties of the semi-group (R^k_τ)_{k∈ℕ} and to Lemma 4.2—in order to increase
 the order of convergence.

$$\begin{aligned} |a_k^{1,1}| &\leq \tau \mathbb{E} \int_{t_k}^{t_{k+1}} |R_\tau(-B)^{1/2+2\kappa}|_{\mathcal{L}(H)} |(-B)^{1-\kappa} R_\tau^k|_{\mathcal{L}(H)} |y|_H |(-B)^{1/2-\kappa} Du(T-t, \tilde{Y}(t))| dt \\ &\leq C |y| \tau \tau^{-1/2-2\kappa} t_k^{-1+\kappa} \int_{t_k}^{t_{k+1}} (1+(T-t)^{-1/2+\kappa}) e^{-\tilde{\mu}(T-t)} \mathbb{E} (1+|\tilde{Y}(t)|^2) dt, \end{aligned}$$

thanks to Lemma 3.2 and Proposition 5.1. By taking expectation, thanks to Lemma 4.3 we have

$$\begin{split} \sum_{k=1}^{m-1} |a_k^{1,1}| &\leq C |y| \tau^{1/2 - 2\kappa} \sum_{k=1}^{m-1} \frac{1}{t_k^{1-\kappa}} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{1/2-\kappa}} \right) e^{-\tilde{\mu}(T-t)} (1+|y|^2) dt \\ &\leq C_\kappa (1+|y|^3) \tau^{1/2 - 2\kappa} \int_0^T \frac{1}{t^{1-\kappa} (T-t)^{1/2-\kappa}} dt \\ &\leq T^{-(1/2 - 2\kappa)} C_\kappa (1+|y|^3) \tau^{1/2 - 2\kappa} \int_0^1 \frac{1}{s^{1-\kappa} (1-s)^{1/2-\kappa}} ds, \end{split}$$

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and we obtain

$$\sum_{k=1}^{m-1} |a_k^{1,1}| \le C(1+|y|^3) T^{-(1/2-2\kappa)} \tau^{1/2-2\kappa}.$$
(31)

2. *Estimate of* $a_k^{1,2}$ First we write

$$\begin{aligned} |a_k^{1,2}| &\leq C\tau \mathbb{E} \int_{t_k}^{t_{k+1}} |R_{\tau}(-B)^{1/2+2\kappa}|_{\mathcal{L}(H)} |\tau(-B)^{1-\kappa} \\ &\times \sum_{l=0}^{k-1} R_{\tau}^{k-l} G(Y_l) ||(-B)^{1/2-\kappa} Du(T-t, \tilde{Y}(t))| dt \end{aligned}$$

Using Lemma 3.2, we can prove the following useful inequality: for $\tau \le \tau_0$ and any $k \ge 1$

$$\tau \sum_{l=1}^{k} \frac{1}{(l\tau)^{1-\kappa}} \frac{1}{(1+\mu_0 \tau)^{l\kappa}} \le C_{\kappa}.$$
(32)

Indeed,

$$\tau \sum_{l=1}^{k} \frac{1}{(l\tau)^{1-\kappa}} \frac{1}{(1+\mu_{0}\tau)^{l\kappa}} \leq C \int_{0}^{t_{k}} \frac{1}{t^{1-\kappa}} \frac{1}{(1+\mu_{0}\tau)^{\kappa}} dt$$
$$\leq \int_{0}^{\infty} \frac{1}{t^{1-\kappa}} e^{-t\frac{\kappa}{\tau} \log(1+\mu_{0}\tau)} dt$$
$$\leq \int_{0}^{\infty} \frac{1}{s^{1-\kappa}} e^{-s} ds \left(\frac{\tau}{\kappa \log(1+\mu_{0}\tau)}\right)^{\kappa}$$
$$\leq C_{\kappa}.$$

Since G is supposed to be bounded, the estimate (32) yields

$$|\tau(-B)^{1-\kappa} \sum_{l=0}^{k-1} R_{\tau}^{k-l} G(Y_l)| \le C \|G\|_{\infty} \tau \sum_{l=1}^{k} \frac{1}{(l\tau)^{1-\kappa}} \frac{1}{(1+\mu_0 \tau)^{l\kappa}} \le C_{\kappa}.$$

With Lemma 4.3 and Proposition 5.2, we can now write

$$|a_k^{1,2}| \le C(1+|y|^2)\tau^{1/2-2\kappa} \int_{t_k}^{t_{k+1}} \left(1+\frac{1}{(T-t)^{1/2-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt$$

and we get

$$\sum_{k=1}^{m-1} |a_k^{1,2}| \le C(1+|y|^2)\tau^{1/2-2\kappa}.$$
(33)

3. *Estimate of* $a_k^{1,3}$

The analysis of this term is more complicated. We recall that since noise is white in space, for any t > 0, $\mathbb{E}|(-B)^{\gamma}W^{B}(t)|^{2} < +\infty$ if and only if $\gamma < 1/4$; as

a consequence, the strategy used above to control $a_k^{1,1}$ can not be used directly otherwise we could only obtain order 1/4.

In [6], an integration by parts formula is used to deal with the lack of regularity of the stochastic integral appearing in the definition of $a_k^{1,3}$. An additionnal difficulty arises in our situation because the estimate given in Lemma 4.5 is not uniform with respect to time. Instead, we remark that the two problems—lack of regularity and bad time dependence—do not occur at the same time; therefore a decomposition of the interval $[0, t_k]$ into $[0, t_k - 1]$ and $[t_k - 1, t_k]$ —for k large enough—can help to treat the problems separately.

Let us explain more concretely the situation at the continuous time level: we have to treat—at the finite dimensional approximation level, but with a bound independent from dimension—an expression involving $B^2 \int_0^t e^{(t-s)B} dW(s)$. The idea to get rid of this expression is to use an integration by parts formula on the whole interval. We can also see that we can do this integration by parts only on a subinterval of size independent from *t*: indeed given $T_0 > 0$, if $t \ge T_0$ we have a decomposition

$$\int_0^t e^{(t-s)B} dW(s) = \int_0^{t-T_0} e^{(t-s)B} dW(s) + \int_{t-T_0}^t e^{(t-s)B} dW(s).$$

The first term is equal to $e^{T_0B} \int_0^{t-T_0} e^{(t-T_0-s)B} dW(s)$, so that thanks to regularization properties of the semi-group (e^{tB}) —see Proposition 2.4—multiplication by B^2 is possible—in other words we do not require an integration by parts to get a bound independent from the dimension; to treat the second term, the lack of regularity remains but can still be treated by the integration by parts, with the advantage of involving a smaller interval—of size T_0 , where at the discrete time level below we can use a uniform control of all quantities, thanks to Lemma 4.5. Then the same idea of "sharing" B^2 can be used again to get order 1/2.

For the discrete-time setting, the choice of T_0 is important: below, we need to use the second estimate of Lemma 3.2 with $\beta = 4 + 1/2 + \kappa$, with $0 < \kappa < 1/2$. Given $\tau_0 > 0$, we then choose $T_0 = 5\tau_0$.

Let us now explain develop this program for the discrete time situation: by using Eq. 10, we make the decomposition

$$\begin{split} a_k^{1,3} &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < R_\tau B^2 \sqrt{\tau} \sum_{l=0}^{k-1} R_\tau^{k-l} \chi_{l+1}, Du(T-t, \tilde{Y}(t)) > dt \\ &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < \int_0^{t_k} R_\tau B^2 R_\tau^{k-l_s} dW(s), Du(T-t, \tilde{Y}(t)) > dt \\ &= -\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < \int_0^{(t_k - T_0) \vee 0} R_\tau B^2 R_\tau^{k-l_s} dW(s), Du(T-t, \tilde{Y}(t)) > dt \\ &- \tau \mathbb{E} \int_{t_k}^{t_{k+1}} < \int_{(t_k - T_0) \vee 0}^{t_k} R_\tau B^2 R_\tau^{k-l_s} dW(s), Du(T-t, \tilde{Y}(t)) > dt. \end{split}$$

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For the first term—which is equal to 0 when $t_k < T_0$ —we use the Cauchy-Schwarz inequality and we directly get

$$\begin{split} |\mathbb{E} < \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau} B^{2} R_{\tau}^{k-l_{s}} dW(s), Du(T-t, \tilde{Y}(t)) > | \\ \leq \left(\mathbb{E} |\int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau} B^{2} R_{\tau}^{k-l_{s}} dW(s)|^{2} \right)^{1/2} (\mathbb{E} |Du(T-t, \tilde{Y}(t))|^{2})^{1/2} \\ \leq C(1+|y|^{2}) e^{-c(T-t)}, \end{split}$$

thanks to Lemmas 5.6 and 4.3, and to the following inequality—we remark that in the integral below $t_{k-l_s} \ge T_0$:

$$\begin{split} \mathbb{E} | \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau} B^{2} R_{\tau}^{k-l_{s}} dW(s) |^{2} \\ &= \int_{0}^{(t_{k}-T_{0})\vee 0} |R_{\tau} B^{2} R_{\tau}^{k-l_{s}} |_{\mathcal{L}_{2}(H)}^{2} ds \\ &= \int_{0}^{(t_{k}-T_{0})\vee 0} \operatorname{Tr}(R_{\tau}^{2} B^{4} R_{\tau}^{2(k-l_{s})}) ds \\ &\leq \int_{0}^{(t_{k}-T_{0})\vee 0} |R_{\tau}^{2+(k-l_{s})}|_{\mathcal{L}(H)} |(-B)^{4+1/2+\kappa} R_{\tau}^{(k-l_{s})}|_{\mathcal{L}(H)} ds \operatorname{Tr}((-B)^{-1/2-\kappa}) \\ &\leq C \int_{0}^{(t_{k}-T_{0})\vee 0} \frac{1}{(1+\mu_{0}\tau)^{k-l_{s}} t_{k-l_{s}}^{4+1/2+\kappa}} ds \\ &\leq C \int_{0}^{(t_{k}-T_{0})\vee 0} \frac{1}{(1+\mu_{0}\tau)^{k-l_{s}}} ds \\ &\leq C \int_{0}^{+\infty} \frac{1}{(1+\mu_{0}\tau)^{s/\tau}} ds \\ &\leq C, \end{split}$$

when $\tau \leq \tau_0$, where we have used the two inequalities of Lemma 3.2. Indeed, $\tau \leq \tau_0$ and $t_{k-l_s} \geq T_0$ imply $k - l_s \geq \frac{T_0}{\tau_0} = 5 \geq 4 + 1/2 + \kappa$. Then

$$\begin{aligned} \left| \tau \mathbb{E} \int_{t_k}^{t_{k+1}} &< \int_0^{(t_k - T_0) \vee 0} R_\tau B^2 R_\tau^{k - l_s} dW(s), Du(T - t, \tilde{Y}(t)) > dt \right| \\ &\leq C \int_{t_k}^{t_{k+1}} e^{-\tilde{\mu}(T - t)} dt (1 + |y|^2) \tau. \end{aligned}$$

For the second term, we use the integration by parts formula of Lemma 4.4 to get

$$\tau \mathbb{E} \int_{t_k}^{t_{k+1}} < \int_{(t_k - T_0) \vee 0}^{t_k} R_{\tau} B^2 R_{\tau}^{k - l_s} dW(s), Du(T - t, \tilde{Y}(t)) > dt$$

= $-\tau \mathbb{E} \int_{t_k}^{t_{k+1}} \int_{(t_k - T_0) \vee 0}^{t_k} \operatorname{Tr} \left(R_{\tau}^{k - l_s} B^2 R_{\tau} D^2 u(T - t, \tilde{Y}(t)) \mathcal{D}_s \tilde{Y}(t) \right) ds dt.$

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If $h \in H$, $(t_k - T_0) \lor 0 \le s \le t_k \le t < t_{k+1}$, with Eq. 13 we see that

$$\mathcal{D}_s^h \tilde{Y}(t) = \mathcal{D}_s^h Y_k + \int_{t_k}^t (B_\tau \mathcal{D}_s^h Y_k + R_\tau G(Y_k) \mathcal{D}_s^h Y_k) d\lambda$$
$$+ R_\tau \mathcal{D}_s^h (W(t) - W(t_k))$$
$$= \mathcal{D}_s^h Y_k + (t - t_k) (B_\tau \mathcal{D}_s^h Y_k + R_\tau G(Y_k) \mathcal{D}_s^h Y_k)$$

Therefore, since $|\tau B_{\tau}|_{\mathcal{L}(H)} \leq C$

$$|\mathcal{D}_s^h \tilde{Y}(t)|_{\beta} \le c |\mathcal{D}_s^h Y_k|_{\beta},$$

and taking supremum over h with $|h| \le 1$ we get

$$|(-B)^{\beta}\mathcal{D}_{s}\tilde{Y}(t)|_{\mathcal{L}(H)} \leq c|(-B)^{\beta}\mathcal{D}_{s}Y_{k}|_{\mathcal{L}(H)}.$$
(34)

The last quantity is estimated thanks to Lemma 4.5:

$$|\mathcal{D}_{s}^{h}Y_{k}|_{\beta} \leq C(1+L_{G}\tau)^{k-l_{s}}\left(1+\frac{1}{(1+\mu_{0}\tau)^{(1-\beta)(k-l_{s})}t_{k-l_{s}}^{\beta}}\right)|h|.$$

When $\tau \leq \tau_0$ and $(t_k - T_0) \vee 0 \leq s \leq t_k \leq t < t_{k+1}$, we see that $(1 + L_G \tau)^{k-l_s}$ is bounded by a constant.

We can then control the second term of $a_k^{1,3}$ with

$$\begin{split} \tau \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{(t_{k}-T_{0})\vee 0}^{t_{k}} |R_{\tau}(-B)^{1/2+2\kappa}|_{\mathcal{L}(H)} |(-B)^{1-3\frac{\kappa}{2}} R_{\tau}^{k-l_{s}}|_{\mathcal{L}(H)} \mathrm{Tr}((-B)^{-1/2-\frac{\kappa}{2}}) \\ &\times |(-B)^{1/2-\kappa/2} D^{2} u(T-t, \tilde{Y}(t))(-B)^{1/2-\kappa/2}|_{\mathcal{L}(H)} |(-B)^{\kappa} \mathcal{D}_{s} \tilde{Y}(t)|_{\mathcal{L}(H)} ds dt \\ &\leq C \tau^{1/2-2\kappa} \int_{t_{k}}^{t_{k+1}} \int_{(t_{k}-T_{0})\vee 0}^{t_{k}} t_{k-l_{s}}^{-1+3\frac{\kappa}{2}} \frac{1}{(1+\mu_{0}\tau)^{(k-l_{s})3\frac{\kappa}{2}}} \\ &\times \left(1+t_{k-l_{s}}^{-\kappa} \frac{1}{(1+\mu_{0}\tau)^{(k-l_{s})(1-\kappa)}}\right) ds \\ &\times \left(1+\frac{1}{(T-t)^{\eta}}+\frac{1}{(T-t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T-t)}(1+|y|^{2}) dt, \end{split}$$

using Proposition 5.2 and Lemmas 4.5 and 3.2. On the one hand, we have

$$\int_{(t_k-T_0)\vee 0}^{t_k} t_{k-l_s}^{-1+3\frac{\kappa}{2}} \frac{1}{(1+\mu_0\tau)^{(k-l_s)3\frac{\kappa}{2}}} ds \leq \int_0^{t_k} \frac{1}{s^{1-3\frac{\kappa}{2}}} \frac{1}{(1+\mu_0\tau)^{3\frac{\kappa}{2}s/\tau}} ds \leq C < +\infty,$$

for $\tau \le \tau_0$, thanks to Eq. 32. On the other hand,

$$\sum_{k=1}^{m-1} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{\eta}} + \frac{1}{(T-t)^{1-\kappa}} \right) e^{-\tilde{\mu}(T-t)} dt$$
$$\leq \int_0^{+\infty} \left(1 + \frac{1}{t^{\eta}} + \frac{1}{t^{1-\kappa}} \right) e^{-\tilde{\mu}t} dt < +\infty.$$

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Therefore

$$\sum_{k=1}^{m-1} |a_k^{1,3}| \le C(1+|y|^2)\tau^{1/2-2\kappa}.$$
(35)

6.2.2 Estimate of a_k^2

We decompose a_k^2 using the definition of \tilde{Y} —see Eq. 13:

$$a_{k}^{2,1} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} (t - t_{k}) < BB_{\tau}Y_{k}, Du(T - t, \tilde{Y}(t)) > dt$$

$$a_{k}^{2,2} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} (t - t_{k}) < BR_{\tau}G(Y_{k}), Du(T - t, \tilde{Y}(t)) > dt$$

$$a_{k}^{2,3} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} < \int_{t_{k}}^{t} BR_{\tau}dW(s), Du(T - t, \tilde{Y}(t)) > dt;$$

then $a_k^2 = a_k^{2,1} + a_k^{2,2} + a_k^{2,3}$.

1. Estimate of $a_k^{2,1}$

Since $BB_{\tau} = R_{\tau}B^2$, $a_k^{2,1}$ is bounded by the same expression as a_k^1 : by Eqs. 31, 33, 35 we have

$$\sum_{k=1}^{m-1} |a_k^{2,1}| \le C(1+|y|^3)(1+T^{-(1/2-2\kappa)})\tau^{1/2-2\kappa}.$$
(36)

2. Estimate of $a_k^{2,2}$ We have

$$\begin{aligned} |a_k^{2,2}| &\leq \tau \mathbb{E} \int_{t_k}^{t_{k+1}} |(-B)^{1/2+\kappa} R_\tau|_{\mathcal{L}(H)} |G(Y_k)| |(-B)^{1/2-\kappa} Du(T-t, \tilde{Y}(t))| dt \\ &\leq \|G\|_{\infty} \tau^{1/2-\kappa} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{1/2-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt. \end{aligned}$$

We then have

$$\sum_{k=1}^{m-1} |a_k^{2,2}| \le C \tau^{1/2-\kappa}.$$
(37)

3. Estimate of $a_k^{2,3}$

We again use the integration by parts formula (14) to rewrite $a_k^{2,3}$:

$$a_k^{2,3} = \mathbb{E} \int_{t_k}^{t_{k+1}} < \int_{t_k}^t BR_\tau dW(s), Du(T-t, \tilde{Y}(t)) > dt$$
$$= \mathbb{E} \int_{t_k}^{t_{k+1}} \int_{t_k}^t \operatorname{Tr}(R_\tau BD^2 u(T-t, \tilde{Y}(t))\mathcal{D}_s \tilde{Y}(t)) ds dt$$

From Eq. 13, for $t_k \le s \le t \le t_{k+1}$ we have $\mathcal{D}_s^h \tilde{Y}(t) = R_\tau h$; as a consequence, we do not need to use the same trick as in the control of $a_k^{1,3}$.

Then we have

$$\begin{split} |a_k^{2,3}| &\leq \mathbb{E} \int_{t_k}^{t_{k+1}} (t-t_k) \mathrm{Tr}(R_{\tau} B D^2 u(T-t, \tilde{Y}(t)) R_{\tau}) dt \\ &\leq c \tau \int_{t_k}^{t_{k+1}} |R_{\tau}(-B)^{1/2+\kappa/2}|_{\mathcal{L}(H)} \mathrm{Tr}((-B)^{-1/2-\kappa/2})|(-B)^{\kappa} R_{\tau}|_{\mathcal{L}(H)} \\ &\times |(-B)^{1/2-\kappa/2} D^2 u(T-t, \tilde{Y}(t))(-B)^{1/2-\kappa/2}|_{\mathcal{L}(H)} dt \\ &\leq c (1+|y|^2) \tau^{1/2-3\kappa/2} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{\eta}} + \frac{1}{(T-t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt. \end{split}$$

Therefore

$$\sum_{k=1}^{m-1} |a_k^{2,3}| \le C(1+|y|^2)\tau^{1/2-3\kappa/2}.$$
(38)

With the previous estimates on a^1 and a^2 , we get

$$\sum_{k=1}^{m-1} |a_k| \le C(1+|y|^3)(1+T^{-(1/2-2\kappa)})\tau^{1/2-2\kappa}.$$
(39)

6.3 Estimate of b_k

We have

$$b_{k} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} \langle G(\tilde{Y}(t)) - R_{\tau}G(Y_{k}), Du(T-t, \tilde{Y}(t)) \rangle dt$$

$$= \mathbb{E} \int_{t_{k}}^{t_{k+1}} \langle (I - R_{\tau})G(Y_{k}), Du(T-t, \tilde{Y}(t)) \rangle dt$$

$$+ \mathbb{E} \int_{t_{k}}^{t_{k+1}} \langle G(\tilde{Y}(t)) - G(Y_{k}), Du(T-t, \tilde{Y}(t)) \rangle dt$$

$$:= b_{k}^{1} + b_{k}^{2}.$$

6.3.1 Estimate of b_k^1

This term is easy to treat: we have

$$\begin{split} |b_{k}^{1}| &\leq \mathbb{E} \int_{t_{k}}^{t_{k+1}} |(-B)^{-1/2+\kappa} (I-R_{\tau})|_{\mathcal{L}(H)} |G(Y_{k})|| (-B)^{1/2-\kappa} Du(T-t, \tilde{Y}(t))| dt \\ &\leq C(1+|y|)^{2}) \tau^{1/2-\kappa} \int_{t_{k}}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{1/2-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt, \end{split}$$

where we have used Proposition 5.1, and the following inequality for $0 \le \beta \le 1$:

$$|(-B)^{-\beta}(I-R_{\tau})|_{\mathcal{L}(H)} \le C_{\beta}\tau^{\beta}.$$
 (40)

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Then we see that

$$\sum_{k=1}^{m-1} |b_k^1| \le C(1+|y|)^2) \tau^{1/2-\kappa}.$$
(41)

6.3.2 Estimate of b_k^2

To estimate $|b_k^2|$, we write the scalar product in coordinates with respect to the orthonormal basis (f_i) , and then we expand the terms thanks to the Itô formula.

If we note $G_i = \langle G, f_i \rangle$ and $\partial_i = \langle D, f_i \rangle$, we have

$$< G(\tilde{Y}(t)) - G(Y_k), Du(T-t, \tilde{Y}(t)) > = \sum_i (G_i(\tilde{Y}(t)) - G_i(Y_k))\partial_i u(T-t, \tilde{Y}(t)).$$

The above sum is finite, because we work with finite dimensional approximations. Itô formula gives for $t_k \le t < t_{k+1}$

$$\begin{split} G_i(\tilde{Y}(t)) - G_i(Y_k) &= \frac{1}{2} \int_{t_k}^t \operatorname{Tr}(R_\tau R_\tau^* D^2 G_i(\tilde{Y}(s))) ds \\ &+ \int_{t_k}^t < B_\tau Y_k, DG_i(\tilde{Y}(s)) > ds \\ &+ \int_{t_k}^t < R_\tau G(Y_k), DG_i(\tilde{Y}(s)) > ds \\ &+ \int_{t_k}^t < DG_i(\tilde{Y}(s)), R_\tau dW(s) > . \end{split}$$

We naturally define $b_k^{2,j}$, for $j \in \{1, 2, 3, 4\}$, and we now control each term.

1. Estimate of $b_k^{2,1}$ By definition, we have

$$b_{k}^{2,1} = \int_{t_{k}}^{t_{k+1}} \mathbb{E}\frac{1}{2} \int_{t_{k}}^{t} \sum_{i} \operatorname{Tr}(R_{\tau}R_{\tau}^{*}D^{2}G_{i}(\tilde{Y}(s))) ds \partial_{i}u(T-t,\tilde{Y}(t)) dt.$$

Using the orthonormal basis $(f_k)_k$ given by Assumption 2.2, and recalling that the sums are finite, we can calculate:

$$\begin{split} \sum_{i} \operatorname{Tr}(R_{\tau}R_{\tau}^{*}D^{2}G_{i}(\tilde{Y}(s)))\partial_{i}u(T-t,\tilde{Y}(t)) \\ &= \sum_{i} \operatorname{Tr}(D^{2}G_{i}(\tilde{Y}(s))R_{\tau}R_{\tau}^{*})\partial_{i}u(T-t,\tilde{Y}(t)) \\ &= \sum_{i} \sum_{j} < D^{2}G_{i}(\tilde{Y}(s))\frac{1}{(1+\mu_{j}\tau)^{2}}f_{j}, f_{j} > \partial_{i}u(T-t,\tilde{Y}(t)) \\ &= \sum_{i} \sum_{j} \frac{1}{(1+\mu_{j}\tau)^{2}}D^{2}G_{i}(\tilde{Y}(s)).(f_{j},f_{j})\partial_{i}u(T-t,\tilde{Y}(t)). \end{split}$$

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Using the Cauchy-Schwarz inequality (where *j* is fixed), we get

$$\begin{split} \left| \sum_{i} D^{2} G_{i}(\tilde{Y}(s)).(f_{j}, f_{j}) \partial_{i} u(T - t, \tilde{Y}(t)) \right| \\ & \leq \left(\sum_{i} \frac{|D^{2} G_{i}(\tilde{Y}(s)).(f_{j}, f_{j})|^{2}}{\mu_{i}^{2\eta}} \right)^{1/2} \left(\sum_{i} \mu_{i}^{2\eta} |\partial_{i} u(T - t, \tilde{Y}(t))|^{2} \right)^{1/2} d_{i} u(T - t, \tilde{Y}(t)) |^{2} d_{i} u(T - t, \tilde{Y}(t)) |^{2}$$

The second factor of this expression is $|(-B)^{\eta}Du(T-t, \tilde{Y}(t))|_{H}$; we control it thanks to Proposition 5.1. The first factor is controlled thanks to Assumption 2.5:

$$\left(\sum_{i} \frac{|D^2 G_i(\tilde{Y}(s)).(f_j, f_j)|^2}{\mu_i^{2\eta}}\right)^{1/2} = |(-B)^{-\eta} D^2 G(\tilde{Y}(s)).(f_j, f_j)|$$

$$\leq C|f_j|_H|f_j|_H \leq C,$$

since $(f_j)_j$ is an orthonormal system. Therefore

$$\begin{split} \left| \sum_{i} \operatorname{Tr}(R_{\tau} R_{\tau}^{*} D^{2} G_{i}(\tilde{Y}(s))) \partial_{i} u(T-t, \tilde{Y}(t)) \right| \\ &\leq C(1+|y|^{2}) \left(1 + \frac{1}{(T-t)^{\eta}} \right) e^{-\tilde{\mu}(T-t)} \sum_{j=0}^{\infty} \frac{1}{(1+\mu_{j}\tau)^{2}} \\ &\leq C(1+|y|^{2}) \left(1 + \frac{1}{(T-t)^{\eta}} \right) e^{-\tilde{\mu}(T-t)} \tau^{-1/2-\kappa} \sum_{j=0}^{\infty} \frac{(\mu_{j}\tau)^{1/2+\kappa}}{(1+\mu_{j}\tau)^{2}} \frac{1}{\mu_{j}^{1/2+\kappa}} \\ &\leq C(1+|y|^{2}) \left(1 + \frac{1}{(T-t)^{\eta}} \right) e^{-\tilde{\mu}(T-t)} \tau^{-1/2-\kappa}. \end{split}$$

Then

$$|b_k^{2,1}| \le C(1+|y|^2)\tau^{1/2-\kappa} \int_{t_k}^{t_{k+1}} \left(1+\frac{1}{(T-t)^{\eta}}\right) e^{-\tilde{\mu}(T-t)} dt,$$

and

$$\sum_{k=1}^{m-1} |b_k^{2,1}| \le C(1+|y|^2)\tau^{1/2-\kappa}.$$
(42)

2. Estimate of $b_k^{2,2}$ Thanks to Eqs. 9 and 10, we have

$$b_{k}^{2,2} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < B_{\tau} R_{\tau}^{k} y + B_{\tau} \tau \sum_{l=0}^{k-1} R_{\tau}^{k-l} G(Y_{l}), DG_{i}(\tilde{Y}(s)) > \partial_{i} u(T-t, \tilde{Y}(t)) ds dt$$
$$+ \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < B_{\tau} \int_{0}^{t_{k}} R_{\tau}^{k-l_{\tau}} dW(r), DG_{i}(\tilde{Y}(s)) > \partial_{i} u(T-t, \tilde{Y}(t)) ds dt$$
$$:= b_{k}^{2,2,1} + b_{k}^{2,2,2}.$$

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(i) For the first term, recalling that $B_{\tau} = BR_{\tau}$ and that G is bounded, we have

$$\begin{split} |b_{k}^{2,2,1}| &= \left| \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} < DG(\tilde{Y}(s)) \cdot \left(B_{\tau} R_{\tau}^{k} y + B_{\tau} \tau \sum_{l=0}^{k-1} R_{\tau}^{k-l} G(Y_{l}) \right), \\ Du(T-t, \tilde{Y}(t)) > dsdt \right| \\ &\leq \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} |(-B)^{\kappa} R_{\tau}|_{\mathcal{L}(H)} (|(-B)^{1-\kappa} R_{\tau}^{k} y| \\ &+ \tau \sum_{l=0}^{k-1} |(-B)^{1-\kappa} R_{\tau}^{k-l}|_{\mathcal{L}(H)} |G(Y_{l})|) \\ &\times |Du(T-t, \tilde{Y}(t))| dsdt \\ &\leq C \tau^{1-\kappa} \int_{t_{k}}^{t_{k+1}} (1+|y|^{2}) e^{-\tilde{\mu}(T-t)} dt \\ &\times \left(t_{k}^{-1+\kappa} |y| + \tau \sum_{l=0}^{k-1} \frac{1}{t_{k-l}^{(1-\kappa)}} \frac{1}{(1+\mu_{0}\tau)^{(k-l)\kappa}} \right) \\ &\leq C \tau^{1-\kappa} (1+|y|^{3}) \left(\frac{1}{t_{k}^{1-\kappa}} + 1 \right) \int_{t_{k}}^{t_{k+1}} e^{-\tilde{\mu}(T-t)} dt, \end{split}$$

if $\tau \leq \tau_0$ —see Eq. 32. Therefore

$$\begin{split} \sum_{k=1}^{m-1} |b_k^{2,2,1}| &\leq C\tau^{1-\kappa} (|y|+1) \sum_{k=1}^{m-1} \left(\frac{1}{t_k^{1-\kappa}}+1\right) \int_{t_k}^{t_{k+1}} e^{-\tilde{\mu}(T-t)} dt \\ &\leq C\tau^{1-\kappa} (1+|y|^3) \int_0^T \left(\frac{1}{t^{1-\kappa}}+1\right) e^{-\tilde{\mu}(T-t)} dt \\ &\leq C\tau^{1-\kappa} (1+|y|^3) \left(\int_0^T \frac{1}{t^{1-\kappa}} \frac{C}{(T-t)^{1/2-\kappa}} dt+1\right) \\ &\leq C\tau^{1-\kappa} (1+|y|^3) \left(T^{-(1/2-2\kappa)} \int_0^1 \frac{1}{s^{1-\kappa}} \frac{C}{(1-s)^{1/2-\kappa}} dt+1\right) \\ &\leq C\tau^{1-\kappa} (1+|y|^3) (T^{-(1/2-2\kappa)}+1). \end{split}$$

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(ii) For the second term, we again use an integration by parts, after a decomposition of the time interval—as in the estimates for $a_k^{1,3}$. We introduce the same parameter $T_0 = 5\tau_0$. First,

$$\begin{split} b_{k}^{2,2,2} &= \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < B_{\tau} \int_{0}^{t_{k}} R_{\tau}^{k-l_{\tau}} dW(r), \\ &DG_{i}(\tilde{Y}(s)) > \partial_{i}u(T-t, \tilde{Y}(t)) ds dt \\ &= \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < B_{\tau} \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau}^{k-l_{\tau}} dW(r) \\ &DG_{i}(\tilde{Y}(s)) > \partial_{i}u(T-t, \tilde{Y}(t)) ds dt \\ &+ \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i,j,m} < B_{\tau} \int_{(t_{k}-T_{0})\vee 0}^{t_{k}} R_{\tau}^{k-l_{\tau}} f_{m}, \\ &f_{j} > d\beta_{m}(r) \partial_{j} G_{i}(\tilde{Y}(s)) \partial_{i}u(T-t, \tilde{Y}(t)) ds dt \\ &=: b_{k}^{2,2,2,1} + b_{k}^{2,2,2,2}. \end{split}$$

For $b_k^{2,2,2,1}$, we can work directly and see that

$$\begin{split} |b_{k}^{2,2,2,1}| &\leq |\mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < B_{\tau} \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau}^{k-l_{\tau}} dW(r), \\ DG_{i}(\tilde{Y}(s)) &> \partial_{i}u(T-t, \tilde{Y}(t)) ds dt | \\ &\leq \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \mathbb{E}| < DG(\tilde{Y}(s)) . B_{\tau} \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau}^{k-l_{\tau}} dW(r), \\ Du(T-t, \tilde{Y}(t)) &> |ds dt \\ &\leq \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \left(\mathbb{E} |B_{\tau} \int_{0}^{(t_{k}-T_{0})\vee 0} R_{\tau}^{k-l_{\tau}} dW(r)|^{2} \right)^{1/2} (\mathbb{E} |Du(T-t, \tilde{Y}(t)) > |^{2})^{1/2} ds dt \\ &\leq C\tau \int_{t_{k}}^{t_{k+1}} e^{-\tilde{\mu}(T-t)} (1+|y|^{2}), \end{split}$$

thanks to Lemmas 4.3, 5.6 and to the following estimate for $\tau \leq \tau_0$

$$\mathbb{E}|B_{\tau}\int_{0}^{(t_{k}-T_{0})\vee 0}R_{\tau}^{k-l_{r}}dW(r)|^{2} \leq \mathbb{E}|B^{2}R_{\tau}\int_{0}^{(t_{k}-T_{0})\vee 0}R_{\tau}^{k-l_{r}}dW(r)|^{2} \leq C,$$

thanks to the estimate proved to control $a_k^{1,3}$.

For $b_k^{2,2,2,2}$, we can write thanks to the Malliavin integration by parts (14) and with the chain rule

$$\begin{split} b_{k}^{2,2,2,2} &= \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \sum_{i,j,m} < B_{\tau} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} R_{\tau}^{k-l_{r}} f_{m}, f_{j} > d\beta_{m}(r)\partial_{j}G_{i}(\tilde{Y}(s))\partial_{i}u \\ &\times (T-t, \tilde{Y}(t))dsdt \\ &= \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,j,m,n} < B_{\tau} R_{\tau}^{k-l_{r}} f_{m}, f_{j} > \partial_{j,n}^{2} G_{i}(\tilde{Y}(s)) \\ &< \mathcal{D}_{r}^{m} \tilde{Y}(s), f_{n} > \partial_{i}u(T-t, \tilde{Y}(t))drdsdt \\ &+ \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,j,m,n} < B_{\tau} R_{\tau}^{k-l_{r}} f_{m}, f_{j} > \partial_{j}G_{i}(\tilde{Y}(s))\partial_{i,n}^{2}u \\ &\times (T-t, \tilde{Y}(t)) < \mathcal{D}_{r}^{m} \tilde{Y}(t), f_{n} > drdsdt \\ &= \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,m} D^{2}G_{i}(\tilde{Y}(s))(B_{\tau} R_{\tau}^{k-l_{r}} f_{m}, \mathcal{D}_{r}^{m} \tilde{Y}(s))\partial_{i}u \\ &\times (T-t, \tilde{Y}(t))drdsdt \\ &+ \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,m} < B_{i}(s, t) B_{\tau} R_{\tau}^{k-l_{r}} f_{m}, \mathcal{D}_{r}^{m} \tilde{Y}(t) > drdsdt \\ &= \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,m} < B_{i}(s, t) B_{\tau} R_{\tau}^{k-l_{r}} f_{m}, \mathcal{D}_{r}^{m} \tilde{Y}(t) > drdsdt \\ &= \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i,m} \operatorname{Tr} \left((\mathcal{D}_{r} \tilde{Y}(s))^{*} D^{2} G_{i}(\tilde{Y}(s)) B_{\tau} R_{\tau}^{k-l_{r}} \right) \partial_{i}u \\ &\times (T-t, \tilde{Y}(t)) drdsdt \\ &+ \mathbb{E} \int_{l_{k}}^{t_{k+1}} \int_{l_{k}}^{t} \int_{(l_{k}-T_{0})\vee 0}^{t_{k}} \sum_{i} \operatorname{Tr} \left((\mathcal{D}_{r} \tilde{Y}(t)^{*} \mathcal{B}_{i}(s, t) B_{\tau} R_{\tau}^{k-l_{r}} \right) drdsdt, \end{split}$$

where we define a linear operator on H by

$$<\mathcal{B}_{i}(s,t)h, k> = \sum_{n=0}^{+\infty}\partial_{i,n}^{2}u(T-t,\tilde{Y}(t)) < k, f_{n} >$$
$$= .$$

We have $\sum_{i} < \mathcal{B}_{i}(s, t)h, k > = D^{2}u(T - t, \tilde{Y}(t)).(DG(\tilde{Y}(s)).h, k)$, and

$$\left|\sum_{i} \mathcal{B}_{i}(s,t)|_{\mathcal{L}(H)} \leq \left| DG(\tilde{Y}(s))|_{\mathcal{L}(H)} |D^{2}u(T-t,\tilde{Y}(t))|_{\mathcal{L}(H)} \right| \right|$$

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so we can write, for $(t_k - T_0) \lor 0 \le r \le t_k$

$$\begin{split} & \left| \sum_{i} \operatorname{Tr} \left((\mathcal{D}_{r} \tilde{Y}(t)^{*} \mathcal{B}_{i}(s, t) B_{\tau} R_{\tau}^{k-l_{r}} \right) \right| \\ & \leq |\mathcal{D}_{r} \tilde{Y}(t)|_{\mathcal{L}(H)} \left| \sum_{i} \mathcal{B}_{i}(s, t) \right|_{\mathcal{L}(H)} |(-B)^{1-3\kappa/2} R_{\tau}^{k-l_{r}}|_{\mathcal{L}(H)} \\ & \times |R_{\tau}(-B)^{1/2+2\kappa}|_{\mathcal{L}(H)} \operatorname{Tr}((-B)^{-1/2-\kappa/2}) \\ & \leq C \tau^{-1/2-2\kappa} t_{k-l_{r}}^{-1+3\kappa/2} \frac{1}{(1+\mu_{0}\tau)^{(k-l_{r})3\kappa/2}} e^{-\tilde{\mu}(T-t)}, \end{split}$$

using Proposition 5.2, Lemma 4.5—since $(1 + L_G \tau)^{k-l_r} \le C$ —Lemma 3.2 and estimate (34).

The other term is a little more complicated, because we are not able to control $D^2G(\tilde{Y}(s))$ in *H*. We proceed as in the estimate of $b_k^{2,1}$, and we directly calculate the trace.

$$\begin{split} &\sum_{i} \operatorname{Tr} \left((\mathcal{D}_{r} \tilde{Y}(s))^{*} D^{2} G_{i}(\tilde{Y}(s)) B_{\tau} R_{\tau}^{k-l_{r}} \right) \partial_{i} u(T-t, \tilde{Y}(t)) \bigg| \\ &\leq |\mathcal{D}_{r} \tilde{Y}(s)|_{\mathcal{L}(H)} \left| \sum_{i} \operatorname{Tr} \left(D^{2} G_{i}(\tilde{Y}(s)) B_{\tau} R_{\tau}^{k-l_{r}} \right) \partial_{i} u(T-t, \tilde{Y}(t)) \right| \\ &\leq |\mathcal{D}_{r} \tilde{Y}(s)|_{\mathcal{L}(H)} \sum_{i,j} \frac{|D^{2} G_{i}(\tilde{Y}(s)).(f_{j}, f_{j})|}{\mu_{i}^{\eta}} \frac{\mu_{j}}{(1+\mu_{j}\tau)^{1+k-l_{r}}} \mu_{i}^{\eta} |\partial_{i} u(T-t, \tilde{Y}(t))| \\ &\leq |\mathcal{D}_{r} \tilde{Y}(s)|_{\mathcal{L}(H)} |(-B)^{\eta} D u(T-t, \tilde{Y}(t))|_{H} \sum_{j} |(-B)^{-\eta} D^{2} G(\tilde{Y}(s)).(f_{j}, f_{j})| \\ &\qquad \times \frac{\mu_{j}}{(1+\mu_{j}\tau)^{1+k-l_{r}}}, \end{split}$$

thanks to the Cauchy-Schwarz inequality.

By using the same analysis as in the estimation of $b_k^{2,1}$, we see that the above expression is bounded by

$$C|\mathcal{D}_r\tilde{Y}(s)|_{\mathcal{L}(H)}|(-B)^{\eta}Du(T-t,\tilde{Y}(t))|_H\sum_j\frac{\mu_j}{(1+\mu_j\tau)^{1+k-l_r}};$$

but the last sum is equal to $\text{Tr}(B_{\tau}R_{\tau}^{k-l_{\tau}})$, so that we see that indeed the two expressions in $b_{k}^{2,2,2}$ are bounded by the same expression.

Therefore

$$\begin{split} |b_{k}^{2,2,2,2}| \\ &\leq \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \int_{(t_{k}-T_{0})\vee 0}^{t_{k}} C\tau^{-1/2-2\kappa} t_{k-l_{r}}^{-1+3\kappa/2} \frac{e^{-\tilde{\mu}(T-t)}}{(1+\mu_{0}\tau)^{(k-l_{r})3\kappa/2}} \\ &\quad \times \left(1+\frac{1}{(T-t)^{\eta}}\right) (1+|y|^{2}) dr ds dt \\ &\leq C(1+|y|^{2})\tau^{1/2-2\kappa} \int_{t_{k}}^{t_{k+1}} \left(1+\frac{1}{(T-t)^{\eta}}\right) e^{-\tilde{\mu}(T-t)} dt \int_{0}^{t_{k}} t_{k-l_{r}}^{-1+3\kappa/2} \frac{1}{(1+\mu_{0}\tau)^{(k-l_{r})3\kappa/2}} dr \\ &\leq C(1+|y|^{2})\tau^{1/2-2\kappa} \int_{t_{k}}^{t_{k+1}} \left(1+\frac{1}{(T-t)^{\eta}}\right) e^{-\tilde{\mu}(T-t)} dt, \end{split}$$

as already proved—see Eq. 32. Now gathering estimates for $b_k^{2,2,2,1}$ and $b_k^{2,2,2,2}$, we obtain

$$\sum_{k=1}^{m-1} |b_k^{2,2,2}| \le C(1+|y|^2)\tau^{1/2-2\kappa}.$$
(43)

3. *Estimate of* $b_k^{2,3}$ We have

$$b_{k}^{2,3} = \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} \sum_{i} < R_{\tau} G(Y_{k}), DG_{i}(\tilde{Y}(s)) > \partial_{i} u(T-t, \tilde{Y}(t)) ds dt$$
$$= \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} < Du(T-t, \tilde{Y}(t)), DG(\tilde{Y}(s)).(R_{\tau} G(Y_{k})) > ds dt.$$

Using that G and DG are bounded, we easily see that

$$|b_k^{2,3}| \le C(1+|y|^2)\tau \int_{t_k}^{t_{k+1}} e^{-\tilde{\mu}(T-t)}dt,$$

and that

$$\sum_{k=1}^{m-1} |b_k^{2,3}| \le C(1+|y|^2)\tau.$$
(44)

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4. *Estimate of* $b_k^{2,4}$ We use the integration by parts formula of Lemma 4.4 to get

$$b_k^{2,4} = \int_{t_k}^{t_{k+1}} \int_{t_k}^t \sum_i < DG_i(\tilde{Y}(s)), R_\tau dW(s) > \partial_i u(T-t, \tilde{Y}(t)) dt$$
$$= \mathbb{E} \int_{t_k}^{t_{k+1}} \int_{t_k}^t \operatorname{Tr} \left((\mathcal{D}_s \tilde{Y}(t))^* D^2 u(T-t, \tilde{Y}(t)) DG(\tilde{Y}(s)) R_\tau \right) ds dt$$
$$= \mathbb{E} \int_{t_k}^{t_{k+1}} \int_{t_k}^t \operatorname{Tr} \left(R_\tau D^2 u(T-t, \tilde{Y}(t)) DG(\tilde{Y}(s)) R_\tau \right) ds dt,$$

using the identity $\mathcal{D}_s^h \tilde{Y}(t) = R_\tau h$ when $t_k \leq s \leq t \leq t_{k+1}$, as in the estimate of $a_k^{2,3}$. Now

$$\begin{split} |b_{k}^{2,4}| &\leq \mathbb{E} \int_{t_{k}}^{t_{k+1}} \int_{t_{k}}^{t} |(R_{\tau}(-B)^{1/2+\kappa}|_{\mathcal{L}(H)} |DG(\tilde{Y}(s))|_{\mathcal{L}(H)} |R_{\tau}|_{\mathcal{L}(H)} \\ &\times |D^{2}u(T-t,\tilde{Y}(t))|_{\mathcal{L}(H)} \mathrm{Tr}((-B)^{-1/2-\kappa}) ds dt \\ &\leq C(1+|y|^{2}) \tau^{1/2-\kappa} \int_{t_{k}}^{t_{k+1}} \left(1+\frac{1}{(T-t)^{\eta}}\right) e^{-\tilde{\mu}(T-t)} dt, \end{split}$$

and

$$\sum_{k=1}^{m-1} |b_k^{2,4}| \le C(1+|y|^2)\tau^{1/2-\kappa}.$$
(45)

6.3.3 Estimate of b_k: conclusion

With Eqs. 41, 42, 43, 44 and 45, we get

$$\sum_{k=1}^{m-1} |b_k| \le C \tau^{1/2 - 2\kappa}.$$
(46)

6.4 Estimate of c_k

We have, using the symmetry of R_{τ} ,

$$\frac{1}{2}I - \frac{1}{2}R_{\tau}R_{\tau}^* = R_{\tau}(I - R_{\tau})^* + \frac{1}{2}(I - R_{\tau})(I - R_{\tau})^*,$$

and

$$c_{k} = \frac{1}{2} \mathbb{E} \int_{t_{k}}^{t_{k+1}} \operatorname{Tr}((I - R_{\tau} R_{\tau}^{*}) D^{2} u(T - t, \tilde{Y}(t))) dt$$

$$= \frac{1}{2} \mathbb{E} \int_{t_{k}}^{t_{k+1}} \operatorname{Tr}((I - R_{\tau})(I - R_{\tau})^{*} D^{2} u(T - t, \tilde{Y}(t))) dt$$

$$+ \mathbb{E} \int_{t_{k}}^{t_{k+1}} \operatorname{Tr}(R_{\tau}(I - R_{\tau})^{*} D^{2} u(T - t, \tilde{Y}(t))) dt$$

$$:= c_{k}^{1} + c_{k}^{2}.$$

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6.4.1 Estimate of c_k^1

We have, using inequality (40)

$$\begin{split} |c_k^1| &\leq \frac{1}{2} \mathbb{E} \int_{t_k}^{t_{k+1}} \operatorname{Tr}((-B)^{-1/2+\kappa} (I - R_{\tau})^2 (-B)^{-1/2+\kappa}) \\ &\times |(-B)^{1/2-\kappa} D^2 u(T - t, \tilde{Y}(t)) (-B)^{1/2-\kappa} |_{\mathcal{L}(H)} dt \\ &\leq C(1 + |y|^2) \int_{t_k}^{t_{k+1}} |(-B)^{-1/2+3\kappa} (I - R_{\tau})|_{\mathcal{L}(H)} |I - R_{\tau}|_{\mathcal{L}(H)} \operatorname{Tr}((-B)^{-1/2-\kappa}) \\ &\times \left(1 + \frac{1}{(T - t)^{\eta}} + \frac{1}{(T - t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T - t)} dt \\ &\leq C(1 + |y|^2) \tau^{1/2-3\kappa} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T - t)^{\eta}} + \frac{1}{(T - t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T - t)} dt. \end{split}$$

Then

$$\sum_{k=1}^{m-1} |c_k^1| \le C(1+|y|^2) \tau^{1/2-3\kappa}.$$
(47)

6.4.2 Estimate of c_k^2

We have, using inequality (40)

$$\begin{split} |c_k^2| &\leq \mathbb{E} \int_{t_k}^{t_{k+1}} \mathrm{Tr}((-B)^{-1/2+\kappa} R_\tau (I-R_\tau) (-B)^{-1/2+\kappa}) \\ &\times |(-B)^{1/2-\kappa} D^2 u (T-t, \tilde{Y}(t)) (-B)^{1/2-\kappa} |_{\mathcal{L}(H)} dt \\ &\leq C (1+|y|^2) \int_{t_k}^{t_{k+1}} |(-B)^{-1/2+\kappa} (I-R_\tau) (-B)^{2\kappa} |_{\mathcal{L}(H)} \mathrm{Tr}((-B)^{-1/2-\kappa}) \\ &\times \left(1 + \frac{1}{(T-t)^{\eta}} + \frac{1}{(T-t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt \\ &\leq C (1+|y|^2) \tau^{1/2-3\kappa} \int_{t_k}^{t_{k+1}} \left(1 + \frac{1}{(T-t)^{\eta}} + \frac{1}{(T-t)^{1-\kappa}}\right) e^{-\tilde{\mu}(T-t)} dt. \end{split}$$

Then

$$\sum_{k=1}^{m-1} |c_k^2| \le C(1+|y|^2) \tau^{1/2-3\kappa}.$$
(48)

6.4.3 Estimate of c_k : conclusion

With Eq. 47 and 48, we get

$$\sum_{k=1}^{m-1} |c_k| \le C(1+|y|^2) \tau^{1/2-3\kappa}.$$
(49)

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6.5 Conclusion

We put together estimates (39), (46), (49) and (30); then passing to the limit with respect to dimension, we get the result.

References

- Brémaud, P.: Markov chains. Gibbs fields, Monte Carlo simulation, and queues. Texts in Applied Mathematics, vol. xviii, p. 444. Springer, New York, NY (1999)
- Brézis, H.: Functional analysis. Theory and applications. (Analyse fonctionnelle. Théorie et applications.). Collection Mathématiques Appliquées pour la Maîtrise, p. 248. Masson, Paris (1994)
- Da Prato, G., Zabczyk, J.: Stochastic equations in infinite dimensions. Encyclopedia of Mathematics and Its Applications, vol. 44, xviii, p. 454. Cambridge etc., Cambridge University Press (1992)
- 4. Da Prato, G., Zabczyk, J.: Ergodicity for infinite dimensional systems. London Mathematical Society Lecture Note Series, 229, xi, p. 339. Cambridge University Press, Cambridge (1996)
- 5. Davie, A.M., Gaines, J.G.: Convergence of numerical schemes for the solution of parabolic stochastic partial differential equations. Math. Comput. **70**(233), 121–134 (2001)
- Debussche, A.: Weak approximation of stochastic partial differential equations: the nonlinear case. Math. Comput. 80(273), 89–117 (2011)
- 7. Debussche, A., Hu, Y., Tessitore, G.: Ergodic BSDEs under weak dissipative assumptions. Stoch. Process. Their Appl. **121**(3), 407–426 (2011)
- Debussche, A., Printems, J.: Weak order for the discretization of the stochastic heat equation. Math. Comput. 78(266), 845–863 (2009)
- Doeblin, W.: Expose de la théorie des chaînes simples constantes de Markoff à un nombre fini d'etats. Rev. Math. Union Interbalkan. 2, 77–105 (1938)
- Gyongy, I.: Lattice approximations for stochastic quasi-linear parabolic partial differential equations driven by space-time white noise. I. Potential Anal. 9(1), 1–25 (1998)
- Gyongy, I.: Lattice approximations for stochastic quasi-linear parabolic partial differential equations driven by space-time white noise. II. Potential Anal. 11(1), 1–37 (1999)
- Gyongy, I., Nualart, D.: Implicit scheme for stochastic parabolic partial differential equations driven by space-time white noise. Potential Anal. 7(4), 725–757 (1997)
- Hausenblas, E.: Approximation for semilinear stochastic evolution equations. Potential Anal. 18(2), 141–186 (2003)
- Kloeden, P.E., Platen, E.: Numerical solution of stochastic differential equations. Applications of Mathematics, vol. 23, xxxv, p. 632. Springer, Berlin (1992)
- Kuksin, S., Shirikyan, A.: A coupling approach to randomly forced nonlinear PDE's. I. Commun. Math. Phys. 221(2), 351–366 (2001)
- Lindvall, E.T.: Lectures on the coupling method. Wiley Series in Probability and Mathematical Statistics. Applied Probability and Statistics, p. 272. Wiley, New York, NY (1992)
- Mattingly, J.C., Stuart, A.M., Higham, D.J.: Ergodicity for SDEs and approximations: locally Lipschitz vector fields and degenerate noise. Stoch. Process. Their Appl. 101(2), 185–232 (2002)
- Mattingly, J.C.: Exponential convergence for the stochastically forced Navier-Stokes equations and other partially dissipative dynamics. Commun. Math. Phys. 230(3), 421–462 (2002)
- Mattingly, J.C., Stuart, A.M., Tretyakov, M.V.: Convergence of numerical time-averaging and stationary measures via Poisson equations. SIAM J. Numer. Anal. 48(2), 552–577 (2010)
- Meyn, S., Tweedie, R.L.: Markov chains and stochastic stability. In: Glynn, P.W. (ed.) Cambridge Mathematical Library, 2nd edn., vol. xviii, p. 594. Cambridge University Press, Cambridge (2009)
- Milstein, G.N.: Numerical integration of stochastic differential equations. Transl. from the Russian. Mathematics and its Applications (Dordrecht), vol. 313, vii, p. 169. Kluwer Academic Publishers, Dordrecht (1994)
- 22. Milstein, G.N., Tretyakov, M.V.: Stochastic numerics for mathematical physics. Scientific Computation, vol. ixx, p. 594. Springer, Berlin (2004)
- Mueller, C.: Coupling and invariant measures for the heat equation with noise. Ann. Probab. 21(4), 2189–2199 (1993)

- Nualart, D.: The Malliavin calculus and related topics, 2nd edn. Probability and Its Applications, vol. xiv, p. 382. Springer, Berlin (2006)
- Printems, J.: On the discretization in time of parabolic stochastic partial differential equations. Monte Carlo Methods Appl. 7(3–4), 359–368 (2001)
- Runst, T., Sickel, W.: Sobolev spaces of fractional order, Nemytskij operators and nonlinear partial differential equations. de Gruyter Series in Nonlinear Analysis and Applications, vol. 3, x, p. 547. de Gruyter, Berlin (1996)
- Sanz-Solé, M.: Malliavin calculus with applications to stochastic partial differential equations. Fundamental Sciences: Mathematics, vol. viii, p. 162. CRC Press, Lausanne: EPFL Press, Boca Raton, FL (2005)
- Talay, D.: Discrétisation d'une équation différentielle stochastique et calcul approché d'espérances de fonctionnelles de la solution. Discretization of a stochastic differential equation and computation of expectations of functions of the solution (1986)
- Talay, D.: Second-order discretization schemes of stochastic differential systems for the computation of the invariant law. Stochastics Stochastics Rep. 29(1), 13–36 (1990)
- Talay, D., Tubaro, L.: Expansion of the global error for numerical schemes solving stochastic differential equations. Stochastic Anal. Appl. 8(4), 483–509 (1990)
- Walsh, J.B.: Finite element methods for parabolic stochastic PDE's. Potential Anal. 23(1), 1–43 (2005)