# WEAK CONVERGENCE FOR A SPATIAL APPROXIMATION OF THE NONLINEAR STOCHASTIC HEAT EQUATION 

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#### Abstract

We find the weak rate of convergence of approximate solutions of the nonlinear stochastic heat equation, when discretized in space by a standard finite element method. Both multiplicative and additive noise is considered under different assumptions.

This extends an earlier result of Debussche in which time discretization is considered for the stochastic heat equation perturbed by white noise. It is known that this equation only has a solution in one space dimension. In order to get results for higher dimensions, colored noise is considered here, besides the white noise case where considerably weaker assumptions on the noise term is needed. Integration by parts in the Malliavin sense is used in the proof. The rate of weak convergence is, as expected, essentially twice the rate of strong convergence.


## 1. Introduction and main result

Let $\mathcal{D} \subset \mathbf{R}^{d}$ be a bounded, convex and polygonal domain. We consider, for $T>0$, the stochastic heat equation with Dirichlet boundary condition, written in abstract form as a stochastic evolution equation in $H=L_{2}(\mathcal{D})$ :

$$
\begin{equation*}
\mathrm{d} X(t)+[A X(t)-f(X(t))] \mathrm{d} t=g(X(t)) \mathrm{d} W(t), t \in(0, T] ; \quad X(0)=X_{0} \tag{1.1}
\end{equation*}
$$

This equation is driven by a cylindrical $Q$-Wiener process $(W(t))_{t \in[0, T]}$ in a filtered probability space $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \in[0, T]}, \mathbf{P}\right)$. The covariance operator $Q$ is selfadjoint and positive semidefinite, not necessarily of finite trace. For technical reasons we consider a deterministic initial value $X_{0} \in H$.

The leading linear operator $A$ is, for simplicity, taken to be $-\Delta$ with domain $\operatorname{dom}(A)=$ $H^{2}(\mathcal{D}) \cap H_{0}^{1}(\mathcal{D})$, where $\Delta=\sum_{k=1}^{d} \partial^{2} / \partial x_{k}^{2}$ is the Laplace operator. It is well known that $-A$ generates an analytic semigroup of bounded linear operators on $H$. We denote it by $(E(t))_{t \geq 0}$. The spaces $\dot{H}^{\beta}=\operatorname{dom}\left(A^{\frac{\beta}{2}}\right)$, defined by fractional powers of $A$, are used to measure the spatial regularity. We denote the norm and inner product in $H=L_{2}(\mathcal{D})$ by $\|\cdot\|$ and $\langle\cdot, \cdot\rangle$.

Let $U, V$ be separable Hilbert spaces and let $\mathcal{L}(U, V)$ denote the Banach space of all linear bounded operators. We denote by $\mathcal{L}_{1}(U, V) \subset \mathcal{L}_{2}(U, V) \subset \mathcal{L}(U, V)$ the subspaces consisting of trace class operators and Hilbert-Schmidt operators, respectively. We use the abbreviations $\mathcal{L}(U)=\mathcal{L}(U, U), \mathcal{L}=\mathcal{L}(H)$ when $H=L_{2}(\mathcal{D})$, and similarly for $\mathcal{L}_{p}, p=1,2$. Central in the theory of stochastic integration is the space $U_{0}=Q^{1 / 2}(H)$. We write $\mathcal{L}_{2}^{0}=\mathcal{L}_{2}\left(U_{0}, H\right)$. By $\mathcal{C}_{\mathrm{b}}^{k}(U, V)$ we denote the space of not necessarily bounded functions from a Banach space $U$ to a Banach space $V$ that have continuous and bounded Fréchet derivatives of orders $1, \ldots, k$. For more precise definitions, see Section 2 below.

We use a "regularity parameter" $\beta$ such that $\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}=\left\|A^{\frac{\beta-1}{2}} Q^{\frac{1}{2}}\right\|_{\mathcal{L}_{2}}<\infty$. If $Q=I$, then $\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}=\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}}<\infty$, if and only if $d=1$ and $\beta<\frac{1}{2}$, see (2.9). We consider two sets of assumptions according to the type of noise term.
A. Additive noise in multiple dimensions. Assume that $f \in \mathcal{C}_{\mathrm{b}}^{2}(H, H), g(x)=I$ for all $x \in H$, and $\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}=\left\|A^{\frac{\beta-1}{2}} Q^{\frac{1}{2}}\right\|_{\mathcal{L}_{2}}<\infty$ for some $\beta \in\left[\frac{1}{2}, 1\right]$.

[^0]B. Multiplicative noise in one dimension. Assume that $f \in \mathcal{C}_{\mathrm{b}}^{2}(H, H), g(x)=B+C x+\tilde{g}(x)$, where $B \in \mathcal{L}, C \in \mathcal{L}(H, \mathcal{L})$, and $\tilde{g} \in \mathcal{C}_{\mathrm{b}}^{2}\left(\dot{H}^{-\frac{1}{2}}, \mathcal{L}\right)$. Moreover, assume that $d=1, Q=I$, and select any $\beta \in\left(0, \frac{1}{2}\right)$.
Under either of these assumptions we have a unique mild solution to (1.1) satisfying the stochastic fixed point equation
\[

$$
\begin{equation*}
X(t)=E(t) X_{0}+\int_{0}^{t} E(t-s) f(X(s)) \mathrm{d} s+\int_{0}^{t} E(t-s) g(X(s)) \mathrm{d} W(s), \quad t \in[0, T] \tag{1.2}
\end{equation*}
$$

\]

One can also show that the solution has spatial regularity of order $\beta$, i.e., it is of the form $X:[0, T] \times \Omega \rightarrow \dot{H}^{\beta}, \mathbf{P}$-almost surely, see Theorem 2.3 below and the discussion preceding it.

In this paper we consider space discretization of equation (1.1) by means of a standard finite element method. Let $\left(S_{h}\right)_{h \in(0,1)}$ be the family of spaces of continuous piecewise linear functions corresponding to a quasi-uniform family of triangulations of $\mathcal{D}$ with $S_{h} \subset H_{0}^{1}(\mathcal{D})$. The parameter $h$ specifies the maximal diameter in the triangulation. Let $P_{h}: H \rightarrow S_{h}$ denote the orthogonal projection. We define the discrete Laplacian as the operator $A_{h}: S_{h} \rightarrow S_{h}$ satisfying the variational equality

$$
\begin{equation*}
\left\langle A_{h} \psi, \chi\right\rangle=\langle\nabla \psi, \nabla \chi\rangle, \quad \forall \psi, \chi \in S_{h} \tag{1.3}
\end{equation*}
$$

The finite element approximation of the elliptic problem $A u=f$ is the unique solution of the equation $A_{h} u_{h}=P_{h} f$. It is known that $\left\|u_{h}-u\right\|=\left\|A_{h}^{-1} P_{h} f-A^{-1} f\right\|=\mathcal{O}\left(h^{2}\right)$ as $h \rightarrow 0$, if $f \in L_{2}(\mathcal{D})$. The semigroup generated by $-A_{h}$ is denoted $\left(E_{h}(t)\right)_{t \geq 0}$. The spatially semidiscrete analogue of (1.1) is to find a process $\left(X_{h}(t)\right)_{t \in(0, T]}$ with values in $S_{h}$ such that

$$
\begin{equation*}
\mathrm{d} X_{h}(t)+\left[A_{h} X_{h}(t)-P_{h} f(X(t))\right] \mathrm{d} t=P_{h} g\left(X_{h}(t)\right) \mathrm{d} W(t), t \in(0, T] ; \quad X_{h}(0)=P_{h} X_{0} \tag{1.4}
\end{equation*}
$$

or in mild form

$$
\begin{align*}
X_{h}(t)= & E_{h}(t) P_{h} X_{0}+\int_{0}^{t} E_{h}(t-s) P_{h} f\left(X_{h}(s)\right) \mathrm{d} s \\
& +\int_{0}^{t} E_{h}(t-s) P_{h} g\left(X_{h}(s)\right) \mathrm{d} W(s), \quad t \in[0, T] \tag{1.5}
\end{align*}
$$

The existence of a unique mild solution can be proved in a similar way as for (1.2). It is also known that we have strong convergence of order $\beta$ under Assumptions A or B, see (2.26). Our goal is to prove weak convergence in the form

$$
\mathbf{E}\left[G(X(T))-G\left(X_{h}(T)\right)\right]=\mathcal{O}\left(h^{2 \beta-\epsilon}\right)
$$

for any $\epsilon>0$ and any testfunction $G \in \mathcal{C}_{\mathrm{b}}^{2}$.
For an exhaustive list of references for approximations of stochastic partial differential equations, see, e.g., [5]. We mention some works related to the situation studied here. Weak convergence of numerical schemes for linear equations with additive noise is treated in [6, [14, [13, and [19]. In the first paper full discretization of the stochastic heat equation is considered for colored noise in multiple dimension, i.e., our Assumption A with $f=0$. Papers [14] and [13] deal with semidiscretization in space and full discretization, respectively, for the linear stochastic heat, Cahn-Hilliard, and wave equations, also with additive colored noise. The fourth paper provides an extension to impulsive noise.

The only results on weak convergence for nonlinear equations are those of [9, [10, [4, [5, [1] and [24]. In the work [9], discretization in time with implicit Euler and Crank-Nicolson schemes is considered for semilinear parabolic equations with additive noise. Paper [10] treats the wave equation with additive white noise, discretized by a leap-frog scheme. This case is a bit different from the others, due to the lack of analyticity of the semigroup for the wave equation in contrast to the heat equation. In [4] semidiscretization in time for the nonlinear stochastic Schrödinger equation with multiplicative white noise is considered.

The papers [4], [6, 14], [13], and [19] express the weak error by means of a Kolmogorov equation after removing the linear term $A X(t)$ by a transformation of variables. This transformation does not work for the nonlinear heat equation. This difficulty is handled in 5 by
means of an integration by parts from the Malliavin calculus. This paper proves weak convergence of temporal semidiscretizations for the nonlinear heat equation with multiplicative noise in one space dimension, i.e, our Assumption B. Under the same assumptions, except for an extra boundedness condition on the nonlinearity, in [1] the method of [5] is exploited to prove weak convergence for the invariant measure of temporally discrete approximations. In [24] the same proof technique is used to study time discretization for the heat equation with additive noise in multiple dimensions, i.e., our Assumption A.

In the present paper we extend the results of [24] and [5] to spatial discretization. Our Assumptions $\mathbf{A}$ and $\mathbf{B}$ coincide with the ones in these two papers, respectively. Therefore we may quote some moment estimates from these papers. One difficulty that arises in connection with the spatial discretization is that the projector $P_{h}$ does not commute with the projector onto eigenspaces of $A$.

In all these works the rate of weak convergence is, up to an arbitrary $\epsilon>0$, twice that of strong convergence. The Malliavin calculus is a useful tool in the study of weak convergence of semilinear equations. It has been utilized in (9, [5, and [15] in completely different ways. It plays a central role in the proof of our Theorem 1.1, following the method of 5]. In the papers [1] and [24] the technique of [5] is also used.

The result of this paper actually concerns the convergence of the law $\mathcal{L}\left(X_{h}(T)\right)=\mathbf{P} \circ$ $\left(X_{h}(T)\right)^{-1}$ of the random variables $\left(X_{h}(T)\right)_{h \in(0,1)}$, as the mesh size parameter $h \rightarrow 0$. We say that the law of $X_{h}(T)$ converges weakly to that of $X(T)$, if $\mathbf{E}\left[G\left(X_{h}(T)\right)\right] \rightarrow \mathbf{E}[G(X(T))]$ as $h \rightarrow 0$, for all test functions $G \in \mathcal{C}_{\mathrm{b}}(H, \mathbf{R})$, the space of all bounded continuous functions on $H$. This convergence follows from the strong convergence $\mathbf{E}\left[\left\|X_{h}(T)-X(T)\right\|^{2}\right]=\mathcal{O}\left(h^{\beta}\right)$, see [16] and the discussion below, and the weak rate obtained is thus $\beta$ under mild assumptions. For $G \in \mathcal{C}_{\mathrm{b}}^{2}(H, \mathbf{R})$, we obtain in this paper the rate of weak convergence $2 \beta-\epsilon$, for an arbitrary $\epsilon>0$.

Theorem 1.1. Assume either Assumption $\boldsymbol{A}$ or Assumption $\boldsymbol{B}$ and let $X$ and $X_{h}$ be the solutions of the equations (1.2) and (1.5), respectively. Then, for every test function $G \in \mathcal{C}_{\mathrm{b}}^{2}(H, \mathbf{R})$ and $\gamma \in[0, \beta)$, we have the convergence

$$
\left|\mathbf{E}\left[G(X(T))-G\left(X_{h}(T)\right)\right]\right|=\mathcal{O}\left(h^{2 \gamma}\right), \quad \text { as } h \rightarrow 0
$$

The weak error is interesting by various reasons. It measures the error made by sampling from an approximate probability law of $X(T)$, rather than the deviation from the trajectory of an exact solution, as for the strong error. The result tells us that the weak error, when approximating the quantity $\mathbf{E}[G(X(T))]$ by $\mathbf{E}\left[G\left(X_{h}(T)\right)\right]$, is decreasing fast as $h \rightarrow 0$ for smooth $G$.

Section 2 is devoted to preliminaries. In Subsection2.1compact operators and tensor products are introduced. We need Schatten classes more general than the trace class and Hilbert-Schmidt operators. In Subsection 2.2 some notation for Fréchet derivatives is fixed. The semigroup framework and basic material on the finite element method are presented in Subsection 2.3. In Subsection 2.4 the Malliavin calculus and stochastic integration is introduced. Subsection 2.5 is about the stochastic equations (1.2) and (1.5). In Section 3 two moment estimates for the Malliavin derivative of $X_{h}(t)$ are proved. Section 4 contains regularity results for the Kolmogorov equation, adapting results from [5] and [24] to our setting. The proof of Theorem 1.1 is given in Section 5

## 2. PRELIMINARIES

2.1. Compact operators and tensor products. Given two separable real Hilbert spaces $\left(U,\langle\cdot, \cdot\rangle_{U}\right)$ and $\left(V,\langle\cdot, \cdot\rangle_{V}\right)$, let $\mathcal{L}(U, V)$ denote the Banach space of all bounded and linear operators $U \rightarrow V$ endowed with the uniform norm. We write $\mathcal{L}(U)=\mathcal{L}(U, U)$. Let $\left(\sigma_{i}\right)_{i \in \mathcal{I}}$ be the collection of singular values of a compact operator $T \in \mathcal{L}(U)$. These are the eigenvalues of the operator $|T|=\left(T T^{*}\right)^{1 / 2}$. The index set $\mathcal{I}$ is finite or countable. Let, for $1 \leq p<\infty$, the

Schatten class $\mathcal{L}_{p}=\mathcal{L}_{p}(U)$ be all $T \in \mathcal{L}(U)$ for which

$$
\begin{equation*}
\|T\|_{\mathcal{L}_{p}}=\left(\sum_{i \in \mathcal{I}} \sigma_{i}^{p}\right)^{\frac{1}{p}}<\infty \tag{2.1}
\end{equation*}
$$

We set by definition $\mathcal{L}_{\infty}=\mathcal{L}$. The Schatten classes are Banach spaces equipped with the norms (2.1). The class $\mathcal{L}_{1}$ is the space of trace class operators. Take an arbitrary ON-basis $\left(e_{n}\right)_{n \in \mathbf{N}} \subset U$. We define the trace of an operator $T \in \mathcal{L}_{1}(U)$ as the quantity

$$
\operatorname{Tr}(T)=\sum_{i \in \mathbf{N}}\left\langle T e_{i}, e_{i}\right\rangle_{U}
$$

It is independent of the particular choice of ON-basis. If $T \in \mathcal{L}_{1}$ and $T \geq 0$, then $\operatorname{Tr}(T)=\|T\|_{\mathcal{L}_{1}}$. In general, the relation

$$
\begin{equation*}
|\operatorname{Tr}(T)| \leq\|T\|_{\mathcal{L}_{1}} \tag{2.2}
\end{equation*}
$$

holds for $T \in \mathcal{L}_{1}$. It follows directly from the definition that $\operatorname{Tr}(T)=\operatorname{Tr}\left(T^{*}\right)$ for $T \in \mathcal{L}_{1}$. Moreover,

$$
\begin{equation*}
\operatorname{Tr}(S T)=\operatorname{Tr}(T S) \tag{2.3}
\end{equation*}
$$

whenever $S \in \mathcal{L}(U, V)$ and $T \in \mathcal{L}(V, U)$ satisfies $S T \in \mathcal{L}_{1}(V)$ and $T S \in \mathcal{L}_{1}(U)$.
More generally, the class $\mathcal{L}_{2}(U, V)$ is the space of Hilbert-Schmidt operators from $U$ to $V$. It is defined as the Hilbert space with the scalar product and norm

$$
\begin{align*}
\langle S, T\rangle_{\mathcal{L}_{2}(U, V)} & =\sum_{i \in \mathbf{N}}\left\langle S e_{i}, T e_{i}\right\rangle_{V}=\operatorname{Tr}\left(T^{*} S\right)=\operatorname{Tr}\left(S T^{*}\right)  \tag{2.4}\\
\|T\|_{\mathcal{L}_{2}(U, V)} & =\left(\sum_{i \in \mathbf{N}}\left\|T e_{i}\right\|_{V}^{2}\right)^{\frac{1}{2}}=\sqrt{\operatorname{Tr}\left(T T^{*}\right)} \tag{2.5}
\end{align*}
$$

The choice of ON-basis $\left(e_{n}\right)_{n \in \mathbf{N}} \subset U$ is arbitrary. For $U=V$ the class $\mathcal{L}_{2}=\mathcal{L}_{2}(U)$ is alone to enjoy this property. For $\mathcal{L}_{p}$ with $p \neq 2$, only an eigenbasis of $|T|$ can be used.

The following Hölder type inequality for Schatten classes holds:

$$
\begin{equation*}
\|S T\|_{\mathcal{L}_{r}} \leq\|S\|_{\mathcal{L}_{p}}\|T\|_{\mathcal{L}_{q}}, \tag{2.6}
\end{equation*}
$$

for $r^{-1}=p^{-1}+q^{-1}, p, q, r \in[1, \infty]$. The border case

$$
\begin{equation*}
\|S T\|_{\mathcal{L}_{r}} \leq\|S\|_{\mathcal{L}}\|T\|_{\mathcal{L}_{r}} \tag{2.7}
\end{equation*}
$$

is included, meaning that $\mathcal{L}_{r}(U)$ is an ideal of the Banach algebra $\mathcal{L}(U)$. Also

$$
\begin{equation*}
\left|\langle S, T\rangle_{\mathcal{L}_{2}}\right|=\left|\operatorname{Tr}\left(S T^{*}\right)\right| \leq\left\|S T^{*}\right\|_{\mathcal{L}_{1}} \leq\|S\|_{\mathcal{L}}\|T\|_{\mathcal{L}_{1}} \tag{2.8}
\end{equation*}
$$

For more about the Schatten classes see [7].
The tensor product space $U \otimes V$ of two Hilbert spaces $U$ and $V$ is a Hilbert space together with a bilinear mapping $U \times V \rightarrow U \otimes V,(u, v) \mapsto u \otimes v$ with dense range and with the inner product

$$
\left\langle u_{1} \otimes v_{1}, u_{2} \otimes v_{2}\right\rangle_{U \otimes V}=\left\langle u_{1}, u_{2}\right\rangle_{U}\left\langle v_{1}, v_{2}\right\rangle_{V}, \quad u_{1}, u_{2} \in U, v_{1}, v_{2} \in V
$$

If $\left(u_{n}\right)_{n \in \mathbf{N}} \subset U$ and $\left(v_{n}\right)_{n \in \mathbf{N}} \subset V$ are ON-bases, then $\left(u_{m} \otimes v_{n}\right)_{m, n \in \mathbf{N}} \subset U \otimes V$ is an ON-basis. The space $U \otimes V$ can be realized in several isomorphic ways. If the tensor product $u \otimes v$ realizes a rank one operator $(u \otimes v) \phi=\langle v, \phi\rangle_{V} u$, for $\phi \in V$, then $U \otimes V \cong \mathcal{L}_{2}(V, U)$. If $U$ and $V$ are spaces of functions of independent variables $x \in \mathcal{D}_{1}$ and $y \in \mathcal{D}_{2}$, then $(u \otimes v)(x, y)=u(x) v(y)$ is also a realization of $U \otimes V$. For instance, if $U=L_{2}(\mathcal{D})$ and $V=L_{2}(\Omega)$, where $\mathcal{D}$ is our spatial domain and $\Omega$ the sample space, then $U \otimes V=L_{2}(\Omega \times \mathcal{D}) \cong L_{2}\left(\Omega, L_{2}(\mathcal{D})\right)$, i.e., $L_{2}(\mathcal{D})$-valued square integrable random variables. For a detailed introduction to tensor products, see [11, Appendix E].
2.2. Fréchet derivatives. Let $\left(U,\|\cdot\|_{U}\right)$ and $\left(V,\|\cdot\|_{V}\right)$ be Banach spaces. By $\mathcal{C}_{\mathrm{b}}^{m}(U, V)$ we denote the space of not necessarily bounded mappings $g: U \rightarrow V$ having $m$ continuous and bounded Fréchet derivatives $D g, D^{2} g, \ldots, D^{m} g$. We endow it with the seminorm $|\cdot|_{\mathcal{C}_{\mathrm{b}}^{m}(U, V)}$, determined as the smallest constant $C \geq 0$ such that

$$
\sup _{x \in U}\left\|D^{m} g(x) \cdot\left(\phi_{1}, \ldots, \phi_{m}\right)\right\|_{V} \leq C\left\|\phi_{1}\right\|_{U} \cdots\left\|\phi_{m}\right\|_{U}, \quad \forall \phi_{1}, \ldots, \phi_{m} \in U
$$

It will be convenient to write $\mathcal{C}_{\mathrm{b}}^{m}=\mathcal{C}_{\mathrm{b}}^{m}(U, V)$. From the context it will be clear what we mean.
Let us consider the important case when $U$ is a Hilbert space and $V=\mathbf{R}$. The Fréchet derivative $D g(x)$ of a function $g: U \rightarrow \mathbf{R}$ is a bounded linear functional on $U$ for fixed $x \in$ $H$ and it can thus be identified by its gradient using the Riesz representation theorem, i.e., $D g(x) \cdot \phi=\langle D g(x), \phi\rangle$. In the same way the second derivative enjoys a representation as a bounded linear operator by the identity $D^{2} g(x) \cdot(\phi, \psi)=\left\langle D^{2} g(x) \phi, \psi\right\rangle$. We will use both representations and it will lead to no confusion.
2.3. The functional analytic framework. We will now introduce the semigroup framework on which our analysis of equations (1.2) and (1.5) relies. Recall from Section 1 that $A=-\Delta$ with $\operatorname{dom}(A)=H^{2}(\mathcal{D}) \cap H_{0}^{1}(\mathcal{D})$ and $H=L_{2}(\mathcal{D})$ with $\mathcal{D} \subset \mathbf{R}^{d}$ a convex polygonal domain. We denote $\|\cdot\|=\|\cdot\|_{H}$ and $\langle\cdot, \cdot\rangle=\langle\cdot, \cdot\rangle_{H}$. The operator $A$ is closed, selfadjoint and positive definite.

There is an orthonormal eigenbasis $\left(\varphi_{i}\right)_{i \in \mathbf{N}} \subset H$ with corresponding eigenvalues $0<\lambda_{1}<$ $\lambda_{2} \leq \cdots \leq \lambda_{i} \rightarrow \infty$, as $i \rightarrow \infty$, for which $A \varphi_{i}=\lambda_{i} \varphi_{i}, i \in \mathbf{N}$. The asymptotics $\lambda_{i} \sim i^{2 / d}$, as $i \rightarrow \infty$, is well known. When the space dimension $d=1$, as in Assumption B, we have

$$
\begin{equation*}
\operatorname{Tr}\left(A^{-\frac{1}{2} \gamma}\right)=\left\|A^{-\frac{1}{2} \gamma}\right\|_{\mathcal{L}_{1}}=\left\|A^{-\frac{1}{4} \gamma}\right\|_{\mathcal{L}_{2}}^{2}<\infty, \quad \forall \gamma>1, \text { if } d=1 \tag{2.9}
\end{equation*}
$$

This means that $\beta \in(0,1)$ under Assumption B.
We define norms of fractional orders by

$$
\|v\|_{\dot{H}^{\beta}}=\left\|A^{\frac{\beta}{2}} v\right\|=\left(\sum_{i \in \mathbf{N}} \lambda_{i}^{\beta}\left\langle v, \varphi_{i}\right\rangle^{2}\right)^{\frac{1}{2}}, \quad \beta \in \mathbf{R}
$$

The spaces $\dot{H}^{\beta}$ are then, for $\beta \geq 0$, defined as $\operatorname{dom}\left(A^{\frac{\beta}{2}}\right)$ and for $\beta<0$ as the closure of $H$ with respect to the $\dot{H}^{\beta}$-norm. The space $\dot{H}^{-\gamma}$ of negative order can be identified with the dual space of $\dot{H}^{\gamma}$. Clearly $\dot{H}^{0}=H$, and it is also well known that $\dot{H}^{1}=H_{0}^{1}(\mathcal{D})$ and $\dot{H}^{2}=H^{2}(\mathcal{D}) \cap H_{0}^{1}(\mathcal{D})$, see [22, Ch. 3].

Let $\left(S_{h}\right)_{h \in(0,1)}$ denote a family of standard finite element spaces of continuous piecewise linear functions corresponding to a quasi-uniform family of triangulations, for which $h$ denotes the largest diameter in the triangulation. Then $S_{h} \subset \dot{H}^{1}$. By $P_{h}$ we denote the orthogonal projector of $H$ onto $S_{h}$. Let $A_{h}: S_{h} \rightarrow S_{h}$ be the unique operator satisfying

$$
\left\langle A_{h} \psi, \chi\right\rangle=\langle\nabla \psi, \nabla \chi\rangle, \quad \forall \psi, \chi \in S_{h}
$$

This is the discrete Laplacian. By definition

$$
\begin{equation*}
\left\|A_{h}^{\frac{1}{2}} \varphi_{h}\right\|=\left\|\nabla \varphi_{h}\right\|=\left\|A^{\frac{1}{2}} \varphi_{h}\right\|=\left\|\varphi_{h}\right\|_{\dot{H}^{1}}, \quad \varphi_{h} \in S_{h} \tag{2.10}
\end{equation*}
$$

Therefore, $P_{h}$ can be extended to $\dot{H}^{-1}$, so that for all $\varphi \in \dot{H}^{-1}$,

$$
\begin{equation*}
\left\|A_{h}^{-\frac{1}{2}} P_{h} \varphi\right\|=\sup _{\psi \in S_{h}} \frac{\langle\varphi, \psi\rangle}{\left\|A_{h}^{\frac{1}{2}} \psi\right\|}=\sup _{\psi \in S_{h}} \frac{\langle\varphi, \psi\rangle}{\left\|A^{\frac{1}{2}} \psi\right\|} \leq \sup _{\psi \in \dot{H}^{1}} \frac{\langle\varphi, \psi\rangle}{\left\|A^{\frac{1}{2}} \psi\right\|}=\left\|A^{-\frac{1}{2}} \varphi\right\| \tag{2.11}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\left\|A_{h}^{\frac{1}{2}} P_{h} \varphi\right\| \leq C\left\|A^{\frac{1}{2}} \varphi\right\|, \quad \varphi \in \dot{H}^{1}, \text { uniformly in } h \tag{2.12}
\end{equation*}
$$

This follows from (2.10) and the well-known fact that $P_{h}$ is bounded with respect to $\|\cdot\|_{\dot{H}^{1}}=$ $\left\|A^{\frac{1}{2}} \cdot\right\|$, when we use a quasi-uniform mesh family. Interpolation between this and (2.11) yields

$$
\begin{equation*}
\left\|A_{h}^{\gamma} P_{h} \varphi\right\| \leq C\left\|A^{\gamma} \varphi\right\|, \quad \varphi \in \dot{H}^{\gamma}, \gamma \in\left[-\frac{1}{2}, \frac{1}{2}\right] \tag{2.13}
\end{equation*}
$$

Furthermore, (2.12) means that $\left\|A_{h}^{\frac{1}{2}} P_{h} A^{-\frac{1}{2}}\right\|_{\mathcal{L}}<\infty$. Hence,

$$
\left\|A^{-\frac{1}{2}} A_{h}^{\frac{1}{2}} P_{h}\right\|_{\mathcal{L}}=\left\|\left(A^{-\frac{1}{2}} A_{h}^{\frac{1}{2}} P_{h}\right)^{*}\right\|_{\mathcal{L}}=\left\|A_{h}^{\frac{1}{2}} P_{h} A^{-\frac{1}{2}}\right\|_{\mathcal{L}} \leq C
$$

so that $\left\|A^{\frac{1}{2}} A_{h}^{-\frac{1}{2}} P_{h} \varphi\right\| \leq C\|\varphi\|$ or

$$
\left\|A^{-\frac{1}{2}} \varphi_{h}\right\| \leq C\left\|A_{h}^{-\frac{1}{2}} \varphi_{h}\right\|, \quad \varphi_{h} \in S_{h} .
$$

Interpolating between this and (2.10) yields

$$
\left\|A^{\gamma} \varphi_{h}\right\| \leq C\left\|A_{h}^{\gamma} \varphi_{h}\right\|, \quad \varphi_{h} \in S_{h}, \gamma \in\left[-\frac{1}{2}, \frac{1}{2}\right] .
$$

Using also (2.13) yields the norm equivalence

$$
\begin{equation*}
c\left\|A_{h}^{\gamma} \varphi_{h}\right\| \leq\left\|A^{\gamma} \varphi_{h}\right\| \leq C\left\|A_{h}^{\gamma} \varphi_{h}\right\|, \quad \varphi_{h} \in S_{h}, \gamma \in\left[-\frac{1}{2}, \frac{1}{2}\right] . \tag{2.14}
\end{equation*}
$$

The interpolations above are valid since $\left(\dot{H}^{\beta}\right)_{\beta \in[-1,1]}$ and $\left(\dot{H}_{h}^{\beta}\right)_{\beta \in[-1,1]}$ are real interpolation spaces, where $\dot{H}_{h}^{\beta}=S_{h}$ with norm $\left\|v_{h}\right\|_{\dot{H}_{h}^{\beta}}=\left\|A_{h}^{\frac{\beta}{2}} v_{h}\right\|$. For positive order this is standard, see for instance [20]. For negative order, let $\beta \in[0,1]$ and notice that

$$
\left[\dot{H}^{0}, \dot{H}^{-1}\right]_{\beta, 2}=\left[\left(\dot{H}^{0}\right)^{*},\left(\dot{H}^{1}\right)^{*}\right]_{\beta, 2}=\left[\dot{H}^{0}, \dot{H}^{1}\right]_{\beta, 2}^{*}=\left(\dot{H}^{\beta}\right)^{*}=\dot{H}^{-\beta} .
$$

We define the Ritz projector $R_{h}: \dot{H}^{1} \rightarrow S_{h}$ to be the orthogonal projection with respect to the $\dot{H}^{1}$-scalar product. Since $\mathcal{D}$ is convex and polygonal it is well known that

$$
\begin{equation*}
\left\|A^{\frac{s}{2}}\left(I-R_{h}\right) A^{-\frac{r}{2}}\right\|_{\mathcal{L}} \leq C h^{r-s}, \quad 0 \leq s \leq 1 \leq r \leq 2 . \tag{2.15}
\end{equation*}
$$

For $P_{h}$ the following error estimate holds

$$
\begin{equation*}
\left\|A^{\frac{s}{2}}\left(I-P_{h}\right) A^{-\frac{r}{2}}\right\|_{\mathcal{L}} \leq C h^{r-s}, \quad 0 \leq s \leq 1, \quad 0 \leq s \leq r \leq 2 . \tag{2.16}
\end{equation*}
$$

For more about the finite element method, see 2] for elliptic equations and 22 for parabolic.
Denote by $N_{h}$ the dimension of $S_{h}$. There is an orthonormal eigenbasis $\left(\varphi_{i}^{h}\right)_{i=1}^{N_{h}} \subset S_{h}$ corresponding to $A_{h}$ with eigenvalues $0<\lambda_{1}^{h} \leq \lambda_{2}^{h} \leq \cdots \leq \lambda_{N_{h}}^{h}$. The operators $-A$ and $-A_{h}$ generate analytic semigroups $(E(t))_{t \geq 0}$ and $\left(E_{h}(t)\right)_{t \geq 0}$, respectively. They are spectrally given by

$$
\begin{equation*}
E(t) v=\sum_{i \in \mathbf{N}} e^{-\lambda_{i} t}\left\langle v, \varphi_{i}\right\rangle \varphi_{i}, \quad v \in H, t \geq 0, \tag{2.17}
\end{equation*}
$$

and

$$
E_{h}(t) v_{h}=\sum_{i=1}^{N_{h}} e^{-\lambda_{i}^{h} t}\left\langle v_{h}, \varphi_{i}^{h}\right\rangle \varphi_{i}^{h}, \quad v_{h} \in S_{h}, t \geq 0 .
$$

The semigroup $\left(E_{h}(t)\right)_{t \geq 0}$ solves the parabolic equation $\dot{u}_{h}+A_{h} u_{h}=0, t \geq 0$, with $u_{h}(0)=P_{h} v$, in the sense that $u_{h}(t)=E_{h}(t) P_{h} v$.

Important for our analysis is the estimate

$$
\begin{equation*}
\left\|A^{\gamma} E(t)\right\|_{\mathcal{L}}+\left\|A_{h}^{\gamma} E_{h}(t) P_{h}\right\|_{\mathcal{L}} \leq C_{\gamma} t^{-\gamma}, \quad \gamma \geq 0, t>0, \quad \text { uniformly in } h . \tag{2.18}
\end{equation*}
$$

It is standard and is enjoyed by all analytic semigroups.
Let $P_{m}$ denote the spectral projection onto the space spanned by the $m$ first eigenvectors $\left(\varphi_{i}\right)_{i=1}^{m}$ of $A$. An easy calculation shows that

$$
\begin{equation*}
\left\|\left(I-P_{m}\right) A^{-r}\right\|_{\mathcal{L}} \leq \lambda_{m}^{-r}, \quad r \geq 0 . \tag{2.19}
\end{equation*}
$$

In our analysis we will use the notation $a \lesssim b$, to mean that there exists a constant $C>0$ such that $a \leq C b$. The constant will never depend on the mesh size $h$.

We will frequently use the following Gronwall lemma:
Lemma 2.1 (Generalized Gronwall lemma). Let $\varphi(t) \geq 0$ be a continuous function on $[0, T]$. If, for some $A, B \geq 0$ and $\alpha, \beta \in[0,1)$, the inequality

$$
\varphi(t) \leq A t^{-\alpha}+B \int_{0}^{t}(t-s)^{-\beta} \varphi(s) \mathrm{d} s
$$

holds, then there is $C=C(B, T, \alpha, \beta)$ such that

$$
\varphi(t) \leq C A t^{-\alpha}, \quad t \in(0, T]
$$

2.4. The stochastic integral and Malliavin calculus. Since we use the Malliavin calculus in the proof of our main result, we outline a framework for the stochastic integral in which this calculus has a natural role. This is an alternative to the more classical procedure, presented in [3]. Our presentation of the Wiener integral relies on [23], and the Malliavin calculus on [8] and [18, where a natural extension of the framework of 21] to Hilbert space valued stochastic integrals using tensor products is presented.

The covariance operator $Q \in \mathcal{L}(H)$ is self adjoint and positive semidefinite. Let $Q^{1 / 2}$ denote the unique positive square root. Let $Q^{-1 / 2}$ be its inverse, restricted to $(\operatorname{ker} Q)^{\perp}$. Define the Hilbert space $U_{0}=Q^{1 / 2}(H)$, equipped with the scalar product $\langle u, v\rangle_{U_{0}}=\left\langle Q^{-1 / 2} u, Q^{-1 / 2} v\right\rangle$. If $\operatorname{Tr}(Q)<\infty$, then the triple $i: U_{0} \hookrightarrow H$ is an abstract Wiener space, where $i$ is the inclusion mapping $i: x \mapsto x$. This triple induces a Gaussian probability measure on $H$ with mean 0 and covariance $Q$. It is referred to as an abstract Wiener measure. The space $U_{0}$ is called the Cameron-Martin space in this context.

Let $I: L_{2}\left([0, T], U_{0}\right) \rightarrow L_{2}(\Omega)$ be an isonormal process: For every $\phi \in L_{2}\left([0, T], U_{0}\right)$ the random variable $I(\phi)$ is centered Gaussian and $I$ has the covariance structure

$$
\mathbf{E}[I(\phi) I(\psi)]=\langle\phi, \psi\rangle_{L_{2}\left([0, T], U_{0}\right)}, \quad \phi, \psi \in L_{2}\left([0, T], U_{0}\right)
$$

The existence of $I$ follows by an application of the Kolmogorov Extension Theorem.
Define, for $u \in U_{0}$, the cylindrical $Q$-Wiener process $W:[0, T] \times U_{0} \rightarrow L_{2}(\Omega)$ by

$$
W(t) u=I\left(\chi_{[0, t]} \otimes u\right), \quad u \in U_{0}, t \in[0, T] .
$$

For $u \in U_{0}$ the process $W(t) u, t \in[0, T]$ is a Brownian motion and given $u, v \in U_{0}$

$$
\mathbf{E}[W(t) u W(s) v]=\min (s, t)\langle u, v\rangle_{U_{0}}
$$

The space of Hilbert-Schmidt operators $\mathcal{L}_{2}^{0}=\mathcal{L}_{2}\left(U_{0}, H\right)$ can be identified with $H \otimes U_{0}$, and $h \otimes u \in \mathcal{L}_{2}^{0}$ for $h \in H, u \in U_{0}$ being the operator $(h \otimes u) v=\langle u, v\rangle_{U_{0}} h, v \in U_{0}$.

We now define the $H$-valued Wiener integral for the simplest possible integrands. Let $\Phi=$ $\chi_{[a, b]} \otimes(h \otimes u) \in L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)$, for $a, b \in[0, T], h \in H$ and $u \in U_{0}$. Then the Wiener integral of $\Phi$ is defined as the $H$-valued random variable

$$
\int_{0}^{T} \Phi(s) \mathrm{d} W(s)=I\left(\chi_{[a, b]} \otimes u\right) \otimes h=(W(b) u-W(a) u) \otimes h \in L_{2}(\Omega, H)
$$

It is not difficult to show that for such integrands the following property, known as Wiener's isometry, holds:

$$
\mathbf{E}\left[\left\|\int_{0}^{T} \Phi(t) \mathrm{d} W(t)\right\|_{H}^{2}\right]=\int_{0}^{T}\|\Phi(t)\|_{\mathcal{L}_{2}^{0}}^{2} \mathrm{~d} t .
$$

The integral extends directly to linear combinations of such integrands by linearity of $I$. By the Wiener isometry, completeness of $L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)$ and classical approximation results for $L_{2}([0, T])$ functions and for compact operators, it extends to all $\Phi \in L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)$.

Let $\mathcal{C}_{\mathrm{p}}^{\infty}\left(\mathbf{R}^{n}\right)$ denote the space of all $\mathcal{C}^{\infty}$-functions over $\mathbf{R}^{n}$ with polynomial growth. We define the family of smooth cylindrical random variables

$$
\mathcal{S}=\left\{X=f\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right): f \in \mathcal{C}_{\mathrm{p}}^{\infty}\left(\mathbf{R}^{N}\right), \phi_{1}, \ldots, \phi_{N} \in L_{2}\left([0, T], U_{0}\right), N \geq 1\right\}
$$

and the corresponding family with values in $H$ as

$$
\mathcal{S}(H)=\left\{F=\sum_{i=1}^{M} X_{i} \otimes h_{i}: X_{1}, \ldots, X_{M} \in \mathcal{S}, h_{1}, \ldots, h_{M} \in H, M \geq 1\right\}
$$

The Malliavin derivative of a random variable in $\mathcal{S}$ with representation $X=f\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right)$ is defined as the $L_{2}\left([0, T], U_{0}\right)$-valued random variable $D X=\sum_{i=1}^{N} \partial_{i} f\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right) \otimes \phi_{i}$. Clearly this is a $U_{0}$-valued stochastic process. We write $D_{t} X=\sum_{i=1}^{N} \partial_{i} f\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right) \otimes$
$\phi_{i}(t)$ for $t \in[0, T]$. The Malliavin derivative of a random variable $F \in \mathcal{S}(H)$ with the representation $F=\sum_{i=1}^{M} f_{i}\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right) \otimes h_{i}$ is given by

$$
D_{t} F=\sum_{i=1}^{M} \sum_{j=1}^{N} \partial_{j} f_{i}\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right) \otimes\left(h_{i} \otimes \phi_{j}(t)\right)
$$

Thus $\left(D_{t} F\right)_{t \in[0, T]}$ is an $\mathcal{L}_{2}^{0}$-valued stochastic process. By $D_{t}^{u} F$ we denote the derivative of $F$ in the direction $u \in U_{0}$ at time $t$, i.e., $D_{t}^{u} F=D_{t} F u$, where

$$
D_{t} F u=\sum_{i=1}^{M} \sum_{j=1}^{N}\left\langle u, \phi_{j}(t)\right\rangle_{U_{0}} \partial_{j} f_{i}\left(I\left(\phi_{1}\right), \ldots, I\left(\phi_{N}\right)\right) \otimes h_{i}
$$

At the very heart of Malliavin calculus is the following integration by parts formula. It tells that, for all $F \in \mathcal{S}(H)$ and $\Phi \in L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)$,

$$
\begin{equation*}
\mathbf{E}\langle D F, \Phi\rangle_{L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)}=\mathbf{E}\left\langle F, \int_{0}^{T} \Phi(t) \mathrm{d} W(t)\right\rangle_{H} \tag{2.20}
\end{equation*}
$$

Thus the Wiener integral is the adjoint of $D: \mathcal{S}(H) \subset L_{2}(\Omega, H) \rightarrow L_{2}\left(\Omega \times[0, T], \mathcal{L}_{2}^{0}\right)$ for deterministic integrands. Formula (2.20) follows from the corresponding formula for real-valued smooth stochastic variables. The derivative operator $D$ is known to be closable. We define the Watanabe Sobolev space $\mathbf{D}^{1,2}(H)$ as the closure of $\mathcal{S}(H)$ with respect to the norm

$$
\|F\|_{\mathbf{D}^{1,2}(H)}=\left(\mathbf{E}\left[\|F\|_{H}^{2}\right]+\mathbf{E}\left[\int_{0}^{T}\left\|D_{t} F\right\|_{\mathcal{L}_{2}^{0}}^{2} \mathrm{~d} t\right]\right)^{\frac{1}{2}}
$$

Denote by $\operatorname{dom}(\delta)$ the elements $\Phi \in L_{2}\left(\Omega \times[0, T], \mathcal{L}_{2}^{0}\right)$ for which $\mathbf{E}\left[\langle D F, \Phi\rangle_{L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)}\right]$ defines a bounded linear functional for $F \in \mathbf{D}^{1,2}(H)$. For any such $\Phi$ the functional $l_{\Phi}(F)=$ $\mathbf{E}\left[\langle D F, \Phi\rangle_{L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)}\right]$ can be extended by continuity to all $F \in L_{2}(\Omega, H)$. The Riesz representation theorem guarantees the existence of an adjoint operator to $D$, namely $\delta: \operatorname{dom}(\delta) \subset$ $L_{2}\left(\Omega \times[0, T], \mathcal{L}_{2}^{0}\right) \rightarrow L_{2}(\Omega, H)$ that satisfies

$$
\begin{equation*}
\mathbf{E}\left[\langle D F, \Phi\rangle_{L_{2}\left([0, T], \mathcal{L}_{2}^{0}\right)}\right]=\mathbf{E}\left[\langle F, \delta(\Phi)\rangle_{H}\right], \quad \forall F \in \mathbf{D}^{1,2}(H) \tag{2.21}
\end{equation*}
$$

This is a natural extension of (2.20) to a much larger class of integrands. In [8, Lemme 2.10] it is proved that for any predictable process $\Phi \in L_{2}\left(\Omega \times[0, T], \mathcal{L}_{2}^{0}\right)$ the action of $\delta$ on $\Phi$ coincides with that of the Itô integral, i.e.,

$$
\delta(\Phi)=\int_{0}^{T} \Phi(t) \mathrm{d} W(t)
$$

Instead of relying on Itô theory we take this as the definition of the Itô integral. We remark that $\operatorname{dom}(\delta)$ contains processes that are not predictable and thus $\delta$ is an extension of the Ito integral to such integrands. In this context $\delta$ is called the Skorohod integral.

The following lemma [5, Lemma 2.1] has a central role in the proof of our main result.
Lemma 2.2. For any random variable $F \in \mathbf{D}^{1,2}(H)$ and any predictable process $\Phi \in L_{2}([0, T] \times$ $\Omega, \mathcal{L}_{2}^{0}$ ) the following integration by parts formula is valid.

$$
\mathbf{E}\left[\left\langle\int_{0}^{T} \Phi(t) \mathrm{d} W(t), F\right\rangle_{H}\right]=\mathbf{E}\left[\int_{0}^{T}\left\langle\Phi(t), D_{t} F\right\rangle_{\mathcal{L}_{2}^{0}} \mathrm{~d} t\right]
$$

Proof. This is just a restatement of (2.21) for $\Phi$ predictable.
A corollary of Lemma 2.2 is the Itô isometry. It reads
(2.22) $\mathbf{E}\left[\left\|\int_{0}^{T} \Phi(t) \mathrm{d} W(t)\right\|_{H}^{2}\right]=\mathbf{E}\left[\int_{0}^{T}\|\Phi(t)\|_{\mathcal{L}_{2}^{0}}^{2} \mathrm{~d} t\right], \quad \forall \Phi \in L_{2}\left([0, T] \times \Omega, \mathcal{L}_{2}^{0}\right)$, predictable.

The Malliavin derivative acts on its adjoint by $D_{s}^{u} \delta(\Phi)=\delta\left(D_{s}^{u} \Phi\right)+\Phi(s) u$ or in terms of the Itô integral $\delta\left(\chi_{[0, t]} \Phi\right)=\int_{0}^{t} \Phi(r) \mathrm{d} W(r)$ with a predictable $\Phi \in L_{2}\left([0, T] \times \Omega, \mathcal{L}_{2}^{0}\right)$ satisfying
$\Phi(t) \in \mathbf{D}^{1,2}\left(\mathcal{L}_{2}^{0}\right)$ for all $t \in[0, T]:$

$$
\begin{equation*}
D_{s}^{u} \int_{0}^{t} \Phi(r) \mathrm{d} W(r)=\int_{0}^{t} D_{s}^{u} \Phi(r) \mathrm{d} W(r)+\Phi(s) u, \quad 0 \leq s \leq t \leq T \tag{2.23}
\end{equation*}
$$

If $s>t$, then $D_{s}^{u} \int_{0}^{t} \Phi(r) \mathrm{d} W(r)=0$, since the integral is $\mathcal{F}_{t}$-measurable. The class of $F \in \mathbf{D}^{1,2}(H)$ that are $\mathcal{F}_{0}$-measurable coincides with the class of constant deterministic random variables. Let $V$ be another separable real Hilbert space and $\sigma \in \mathcal{C}_{\mathrm{b}}^{1}(H, V)$. Then $\sigma(F) \in \mathbf{D}^{1,2}(V)$ and

$$
\begin{array}{rlrl}
D_{t}^{u}(\sigma(F)) & =D \sigma(F) \cdot D_{t}^{u} F, & & u \in U_{0}, F \in \mathbf{D}^{1,2}(H) \\
D_{t}(\sigma(F)) & =D \sigma(F) D_{t} F, & F \in \mathbf{D}^{1,2}(H) \tag{2.25}
\end{array}
$$

2.5. Existence and uniqueness. Existence and uniqueness of a solution to (1.2), under Assumption A with $\beta=1$, is stated as [3, Theorem 7.4]. This is the case when $\operatorname{Tr}(Q)<\infty$. The extension to $\beta \in\left[\frac{1}{2}, 1\right)$ is straight-forward. For Assumption $\mathbf{B}$ existence and uniqueness is given as [3, Theorem 7.6]. By using the methods of [12] and [17] one can show that the regularity in space is of order $\beta$, i.e., the solution $X$ is of the form $[0, T] \times \Omega \rightarrow \dot{H}^{\beta}$, P-a.s.. Recall here that $\beta$ is some number $\beta \in\left[\frac{1}{2}, 1\right]$ under Assumption $\mathbf{A}$ and any number $\beta \in\left(0, \frac{1}{2}\right)$ under Assumption B. The family $\left(X_{h}\right)_{h \in(0,1)}$ of solution processes of the discrete equation (1.5), corresponding to the family of triangulations, is treated analogously and clearly $X_{h}(t) \in S_{h} \subset \dot{H}^{1}, \mathbf{P}$-a.s.. The estimate $\mathbf{E}\left\|A^{\frac{\gamma}{2}} X_{h}(t)\right\|^{2} \leq C\left(1+\left\|X_{0}\right\|^{2}\right), \gamma \in[0,1]$ is uniform in $h$, only for $\gamma \in[0, \beta]$. The strong convergence

$$
\begin{equation*}
\left(\mathbf{E}\left\|X(T)-X_{h}(T)\right\|^{2}\right)^{\frac{1}{2}} \leq C h^{\beta} \tag{2.26}
\end{equation*}
$$

is proved in [16] under the assumption of trace class noise. The proof is easier under Assumptions $\mathbf{A}$ and $\mathbf{B}$. We formulate a qualitative bound for the solution processes in the following theorem.
Theorem 2.3. Under either Assumption $\boldsymbol{A}$ or Assumption $\boldsymbol{B}$ there exists unique predictable solutions $X \in C\left([0, T], L_{2}(\Omega, H)\right)$ and $X_{h} \in C\left([0, T], L_{2}\left(\Omega, S_{h}\right)\right)$ to equation (1.2) and (1.5) respectively. We refer to these solutions as the unique mild solutions of (1.1) and (1.5). There exists a constant $C$, such that the following moment estimate holds

$$
\begin{equation*}
\sup _{t \in[0, T]} \mathbf{E}\|X(t)\|^{2}+\sup _{t \in[0, T]} \mathbf{E}\left\|X_{h}(t)\right\|^{2} \leq C\left(1+\left\|X_{0}\right\|^{2}\right) \tag{2.27}
\end{equation*}
$$

## 3. Estimates of the Malliavin derivative of the solution

We consider the Malliavin derivative of the discrete solution process and prove some estimates needed later. Differentiating the equation (1.5) formally in direction $u \in U_{0}$, using (2.23), (2.24), and the fact that we have a deterministic initial value, yields

$$
\begin{align*}
D_{s}^{u} X_{h}(t)= & E_{h}(t-s) P_{h} g\left(X_{h}(s)\right) u+\int_{s}^{t} E_{h}(t-r) P_{h} D f\left(X_{h}(r)\right) \cdot D_{s}^{u} X_{h}(r) \mathrm{d} r \\
& +\int_{s}^{t} E_{h}(t-r) P_{h}\left(D g\left(X_{h}(r)\right) \cdot D_{s}^{u} X_{h}(r)\right) \mathrm{d} W(r), \quad 0 \leq s \leq t \leq T \tag{3.1}
\end{align*}
$$

This equation is treated much like (1.5) itself. It has a unique solution.
Before we proceed to the estimate of the Malliavin derivative we notice that, by the linear growth of $f$ and $g$, implied by their bounded first derivative, and the moment estimate (2.27) for $X$ and $X_{h}$ yields

$$
\begin{equation*}
\sup _{t \in[0, T]} \mathbf{E}\|f(Y(t))\|^{2}+\sup _{t \in[0, T]} \mathbf{E}\left\|A^{\frac{\beta-1}{2}} g(Y(t))\right\|_{\mathcal{L}_{2}^{0}}^{2} \lesssim 1+\left\|X_{0}\right\|^{2}, \quad Y=X \text { or } X_{h} \tag{3.2}
\end{equation*}
$$

Lemma 3.1. Consider equation (1.5) under Assumption A. Then the Malliavin derivative of $X_{h}$, given as the solution $D_{s} X_{h}$ to equation (3.1), satisfies for some constant $C=C(T)>0$ the bound:

$$
\mathbf{E}\left[\left\|A_{h}^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}}^{2}\right] \leq C, \quad 0 \leq s \leq t \leq T
$$

Proof. We make use of equation (3.1) with $g(x)=I, D g(x)=0$, for the proof and recall that $\beta-1 \in\left[-\frac{1}{2}, 0\right]$. Fix $u \in U_{0}$. Thanks to the Cauchy-Schwarz inequality we get that

$$
\mathbf{E}\left\|D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim\left\|E_{h}(t-s) A_{h}^{\frac{1-\beta}{2}} A_{h}^{\frac{\beta-1}{2}} P_{h} u\right\|^{2}+\int_{s}^{t} \mathbf{E}\left\|E_{h}(t-r) P_{h} D f\left(X_{h}(r)\right) \cdot D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r
$$

In view of (2.13) and the boundedness of $D f$ and $E_{h}(t)$ we have

$$
\begin{equation*}
\mathbf{E}\left\|D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim\left\|A_{h}^{\frac{1-\beta}{2}} E_{h}(t-s) P_{h}\right\|_{\mathcal{L}}^{2}\left\|A^{\frac{\beta-1}{2}} u\right\|^{2}+\int_{s}^{t}|f|_{\mathcal{C}_{\mathrm{b}}^{1}}^{2} \mathbf{E}\left\|D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r \tag{3.3}
\end{equation*}
$$

The analyticity of the semigroup (2.18) yields

$$
\mathbf{E}\left\|D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim(t-s)^{\beta-1}\left\|A^{\frac{\beta-1}{2}} u\right\|^{2}+\int_{s}^{t} \mathbf{E}\left\|D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r
$$

and applying Gronwall's Lemma 2.1, for fixed $s \in[0, t)$, gives

$$
\begin{equation*}
\mathbf{E}\left\|D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim(t-s)^{\beta-1}\left\|A^{\frac{\beta-1}{2}} u\right\|^{2} \tag{3.4}
\end{equation*}
$$

Proceeding as in the proof of (3.3), we obtain also

$$
\mathbf{E}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim\left\|A^{\frac{\beta-1}{2}} u\right\|^{2}+\int_{s}^{t} \mathbf{E}\left\|D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r
$$

Estimate (3.4) is applicable. Thus

$$
\int_{s}^{t} \mathbf{E}\left\|D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r \lesssim \int_{s}^{t}(r-s)^{\beta-1} \mathrm{~d} r\left\|A^{\frac{\beta-1}{2}} u\right\|^{2} \lesssim(t-s)^{\beta}\left\|A^{\frac{\beta-1}{2}} u\right\|^{2}
$$

and hence

$$
\begin{equation*}
\mathbf{E}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim\left\|A^{\frac{\beta-1}{2}} u\right\|^{2} \tag{3.5}
\end{equation*}
$$

Notice that this is uniform with respect to $u \in U_{0}$. We take an ON-basis $\left(u_{i}\right)_{i \in \mathbf{N}} \subset U_{0}$ and compute the $\mathcal{L}_{2}^{0}$-norm according to (2.4). Using Tonelli's Theorem and (3.5) we get that

$$
\begin{aligned}
\mathbf{E}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}}^{2} & =\mathbf{E} \sum_{i \in \mathbf{N}}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s}^{u_{i}} X_{h}(t)\right\|^{2}=\sum_{i \in \mathbf{N}} \mathbf{E}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s}^{u_{i}} X_{h}(t)\right\|^{2} \\
& \lesssim \sum_{i \in \mathbf{N}}\left\|A^{\frac{\beta-1}{2}} u_{i}\right\|^{2}=\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2} .
\end{aligned}
$$

For the white noise case we will need the following lemma that is a space discrete analogue of [5, Lemma 4.3]. Recall that in this case $Q=I, U_{0}=H, \mathcal{L}_{2}^{0}=\mathcal{L}_{2}$

Lemma 3.2. Consider equation (1.5) under Assumption B. Then, for $\gamma \in\left[0, \frac{1}{2}\right)$, the Malliavin derivative satisfies the following estimate:

$$
\mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}}^{2} \leq C(t-s)^{-\gamma}, \quad 0 \leq s \leq t \leq T
$$

Proof. Let $u \in H$, and take norms in (3.1) using the Cauchy-Schwarz inequality and the Itô isometry (2.22) to get

$$
\begin{aligned}
\mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim & \mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} E_{h}(t-s) P_{h} g\left(X_{h}(s)\right) u\right\|^{2} \\
& +\int_{s}^{t} \mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} E_{h}(t-s) P_{h} D f\left(X_{h}(s)\right) \cdot D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r \\
& +\int_{s}^{t} \mathbf{E}\left\|A_{h}^{\frac{1}{4}+\epsilon} A_{h}^{\frac{\gamma}{2}} E_{h}(t-s) A_{h}^{-\frac{1}{4}-\epsilon} P_{h} D g\left(X_{h}(s)\right) \cdot D_{s}^{u} X_{h}(r)\right\|_{\mathcal{L}_{2}}^{2} \mathrm{~d} r
\end{aligned}
$$

For $\epsilon>0$ small enough we have by (2.7) and (2.18)

$$
\begin{aligned}
\mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim & (t-s)^{-\gamma} \sup _{s \in[0, T]} \mathbf{E}\left\|g\left(X_{h}(s)\right)\right\|_{\mathcal{L}}^{2}\|u\|^{2} \\
& +\int_{s}^{t}(t-r)^{-\gamma}|f|_{\mathcal{C}_{\mathrm{b}}^{1}}^{2} \mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r \\
& +\int_{s}^{t}(t-r)^{-\gamma-\frac{1}{2}-2 \epsilon}\left\|A_{h}^{-\frac{1}{4}-\epsilon} P_{h}\right\|_{\mathcal{L}_{2}}^{2}|g|_{\mathcal{C}_{\mathrm{b}}^{1}}^{2} \mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s}^{u} X_{h}(r)\right\|^{2} \mathrm{~d} r .
\end{aligned}
$$

By (2.9), (2.7) and (2.13) we have

$$
\left\|A_{h}^{-\frac{1}{4}-\epsilon} P_{h}\right\|_{\mathcal{L}_{2}} \lesssim\left\|A_{h}^{-\frac{1}{4}-\epsilon} P_{h} A^{\frac{1}{4}+\epsilon}\right\|_{\mathcal{L}}\left\|A^{-\frac{1}{4}-\epsilon}\right\|_{\mathcal{L}_{2}} \lesssim\left\|A^{-\frac{1}{4}-\epsilon}\right\|_{\mathcal{L}_{2}}<\infty
$$

and by Gronwall's Lemma 2.1 and (3.2)

$$
\mathbf{E}\left\|A_{h}^{\frac{\gamma}{2}} D_{s}^{u} X_{h}(t)\right\|^{2} \lesssim(t-s)^{-\gamma}\left(1+\left\|X_{0}\right\|^{2}\right)\|u\|^{2}
$$

## 4. Regularity results for the Kolmogorov equation

In [5], 1] and [24, in the case of discretization in time, the proofs of the weak convergence is performed for finite-dimensional spectral Galerkin approximations. The use of the Itô formula and the Kolmogorov equation is in this way justified. The estimates are uniform in the dimension $m \in \mathbf{N}$ of the approximation space and they thus hold in the limit. The approximation is not made explicit in the proof. For the discretization in space some more care need to be taken. This is due to the fact that the operators $A$ and $A_{h}$ do not commute.

Recall that $P_{m}$ is the projection onto the subspace of $H_{m} \subset H$ spanned by the first $m \in \mathbf{N}$ eigenvectors $\left(\varphi_{i}\right)_{i=1}^{m}$ of $A$. Let $A_{m}=P_{m} A P_{m}=A P_{m}=P_{m} A$. By $\left(E_{m}(t)\right)_{t \geq 0}$ we denote the semigroup generated by $-A_{m}$, i.e., it is given by the $m$ first terms in the spectral representation (2.17) of $(E(t))_{t \geq 0}$.

We denote by $X_{m}^{x}$ the solution of equation
$X_{m}^{x}(t)=E_{m}(t) P_{m} x+\int_{0}^{t} E_{m}(t-s) P_{m} f\left(X_{m}^{x}(s)\right) \mathrm{d} s+\int_{0}^{t} E_{m}(t-s) P_{m} g\left(X_{m}^{x}(s)\right) \mathrm{d} W(s), t \in[0, T]$.
Define the function $u_{m}(t, x)=\mathbf{E}\left[G\left(X_{m}^{x}(t)\right)\right]$ for $t \in[0, T]$ and $x \in H$. Note that $u\left(t, P_{m} x\right)=$ $u(t, x)$ for $x \in H$. It is well known, see e.g. [3, Theorem 9.16], that $u_{m}:[0, T] \times H \rightarrow \mathbf{R}$ is a solution to the Kolmogorov equation

$$
\begin{array}{ll}
\dot{u}_{m}(t, x)+L_{m} u_{m}(t, x)=0, & (t, x) \in(0, T] \times H \\
u_{m}(0, x)=G\left(P_{m} x\right), & x \in H,
\end{array}
$$

where the Markov generator $L_{m}$ is given by

$$
\left(L_{m} v\right)(x)=\left\langle A_{m} x-P_{m} f(x), D v(x)\right\rangle-\frac{1}{2} \operatorname{Tr}\left(P_{m} g(x) Q g^{*}(x) P_{m} D^{2} v(x)\right), \quad x \in H
$$

The proof of Theorem 1.1 relies heavily on estimates of the derivatives $D u_{m}$ and $D^{2} u_{m}$ of $u_{m}$ of the form: for some $\alpha>0$ we have

$$
\begin{align*}
& \sup _{x \in H}\left\|A^{\lambda} D u_{m}(t, x)\right\| \leq C t^{-\lambda}|G|_{C_{\mathrm{b}}^{1}}, \quad t \in(0, T], \lambda \in[0, \alpha)  \tag{4.1}\\
& \sup _{x \in H}\left\|A^{\lambda} D^{2} u_{m}(t, x) A^{\rho}\right\|_{\mathcal{L}} \leq C t^{-(\rho+\lambda)}|G|_{C_{\mathrm{b}}^{2}}, \quad t \in(0, T], \lambda, \rho \in[0, \alpha), \lambda+\rho<1 . \tag{4.2}
\end{align*}
$$

In the case of colored noise it turns out that we need $\alpha \geq(1+\beta) / 2$ to obtain convergence of the right rate. So far, to our knowledge, there is no satisfactory result in this direction for multiplicative noise. But for additive colored noise, case $\mathbf{A}$, the situation is much easier and the estimates hold for $\alpha=1$, see Lemma 3.3 in [24]. For the white noise case the estimates are stated as Lemma 4.4 and Lemma 4.5 in [5] with $\alpha=\frac{1}{2}$. Thus, in case $\mathbf{A}$ we have $\beta \in\left[\frac{1}{2}, 1\right]$, (4.1) and (4.2) with $\alpha=1$, and in case $\mathbf{B}$ we have $\beta \in\left[0, \frac{1}{2}\right)$ and (4.1) and (4.2) with $\alpha=\frac{1}{2}$.

Since we have the operator $A$ in (4.1) and (4.2) instead of the more natural choice $A_{m}$ we outline their proofs. We will use that $D u(t, x) \cdot \phi=\mathbf{E}\left[D G\left(X_{m}^{x}(t)\right) \cdot \eta_{m}^{\phi, x}(t)\right]$, where

$$
\begin{aligned}
\eta_{m}^{\phi, x}(t)= & E_{m}(t) P_{m} \phi+\int_{0}^{t} E_{m}(t-s) P_{m} D f\left(X_{m}^{x}(s)\right) \cdot \eta_{m}^{\phi, x}(s) \mathrm{d} s \\
& +\int_{0}^{t} E_{m}(t-s) P_{m}\left(D g\left(X_{m}^{x}(s)\right) \cdot \eta_{m}^{\phi, x}(s)\right) \mathrm{d} W(s)
\end{aligned}
$$

In the proofs of Lemma 3.3 in [24] for case $\mathbf{A}$ with $\alpha=1$ and Lemma 4.4 in [5] for the case $\mathbf{B}$ with $\alpha=\frac{1}{2}$ it is proved that

$$
\begin{equation*}
\left(\sup _{x \in H} \mathbf{E}\left\|\eta_{m}^{\phi, x}\right\|^{2}\right)^{\frac{1}{2}} \lesssim t^{-\lambda}\left\|A_{m}^{-\lambda} P_{m} \phi\right\|, \quad t \in(0, T], \lambda \in[0, \alpha) \tag{4.3}
\end{equation*}
$$

Therefore

$$
\begin{aligned}
\left\langle A^{\lambda} D u_{m}(t, x), \psi\right\rangle & =\left\langle D u_{m}(t, x), A^{\lambda} \psi\right\rangle=\mathbf{E}\left[D G\left(X_{m}^{x}(t)\right) \cdot \eta_{m}^{A^{\lambda} \psi, x}(t)\right] \leq|G|_{\mathcal{C}_{\mathrm{b}}^{1}}\left(\mathbf{E}\left\|\eta_{m}^{A^{\lambda} \psi, x}(t)\right\|^{2}\right)^{\frac{1}{2}} \\
& \lesssim|G|_{\mathcal{C}_{\mathrm{b}}^{1}} t^{-\lambda}\left\|A_{m}^{-\lambda} P_{m} A^{\lambda} \psi\right\|=|G|_{\mathcal{C}_{\mathrm{b}}^{1}} t^{-\lambda}\left\|P_{m} \psi\right\| \leq|G|_{\mathcal{C}_{\mathrm{b}}} t^{-\lambda}\|\psi\|
\end{aligned}
$$

implying (4.1).
For (4.2) we notice that

$$
\begin{equation*}
D^{2} u_{m}(t, x) \cdot(\phi, \psi)=\mathbf{E}\left[D^{2} G\left(X_{m}^{x}(t)\right) \cdot\left(\eta_{m}^{\phi, x}(t), \eta_{m}^{\psi, x}(t)\right)+D G\left(X_{m}^{x}(t)\right) \cdot \zeta_{m}^{\phi, \psi, x}(t)\right] \tag{4.4}
\end{equation*}
$$

where

$$
\begin{aligned}
\zeta_{m}^{\phi, \psi, x}(t) & =\int_{0}^{t} E_{m}(t-s) P_{m}\left(D^{2} f\left(X_{m}^{x}(s)\right) \cdot\left(\eta_{m}^{\phi, x}(s), \eta_{m}^{\psi, x}(s)\right)+D f\left(X_{m}^{x}(s)\right) \cdot \zeta_{m}^{\phi, \psi, x}(s)\right) \mathrm{d} s \\
& +\int_{0}^{t} E_{m}(t-s) P_{m}\left(D^{2} g\left(X_{m}^{x}(s)\right) \cdot\left(\eta_{m}^{\phi, x}(s), \eta_{m}^{\psi, x}(s)\right)+D g\left(X_{m}^{x}(s)\right) \cdot \zeta_{m}^{\phi, \psi, x}(s)\right) \mathrm{d} W(s)
\end{aligned}
$$

In the proof of Lemma 3.3 in [24] for case $\mathbf{A}$ with $\alpha=1$ and Lemma 4.5 in [5] for the case $\mathbf{B}$ with $\alpha=\frac{1}{2}$ it is shown that

$$
\begin{equation*}
\left(\sup _{t \in[0, T]} \sup _{x \in H} \mathbf{E}\left\|\zeta_{m}^{\phi, \psi, x}(t)\right\|^{2}\right)^{\frac{1}{2}} \lesssim\left\|A_{m}^{-\rho} P_{m} \phi\right\|\left\|A_{m}^{-\lambda} P_{m} \psi\right\|, \quad \lambda, \rho \in[0, \alpha), \lambda+\rho<1 \tag{4.5}
\end{equation*}
$$

Since $D^{2} u_{m} \cdot(\phi, \psi)=\left\langle D^{2} u_{m} \phi, \psi\right\rangle$ and by (4.4) and the Cauchy-Schwarz inequality

$$
\begin{aligned}
& \left\langle A^{\lambda} D^{2} u_{m}(t, x) A^{\rho} \phi, \psi\right\rangle=\left\langle D^{2} u_{m}(t, x) A^{\rho} \phi, A^{\lambda} \psi\right\rangle \\
& \quad=\mathbf{E}\left[D^{2} G\left(X_{m}^{x}(t)\right) \cdot\left(\eta_{m}^{A^{\lambda} \psi, x}(t), \eta_{m}^{A^{\rho} \phi, x}(t)\right)+D G\left(X_{m}^{x}(t)\right) \cdot \zeta_{m}^{A^{\lambda} \psi, A^{\rho} \phi, x}(t)\right] \\
& \quad \leq|G|_{\mathcal{C}_{\mathrm{b}}^{2}}\left(\mathbf{E}\left\|\eta_{m}^{A^{\lambda} \psi, x}(t)\right\|^{2}\right)^{\frac{1}{2}}\left(\mathbf{E}\left\|\eta_{m}^{A^{\rho} \psi, x}(t)\right\|^{2}\right)^{\frac{1}{2}}+|G|_{\mathcal{C}_{\mathrm{b}}^{1}}\left(\mathbf{E}\left\|\zeta_{m}^{A^{\lambda} \psi, A^{\rho} \phi, x}(t)\right\|^{2}\right)^{\frac{1}{2}}
\end{aligned}
$$

Applying (4.3) and (4.5) yields

$$
\begin{aligned}
\left\langle A^{\lambda} D^{2} u_{m}(t, x) A^{\rho} \phi, \psi\right\rangle & \lesssim\left(|G|_{\mathcal{C}_{\mathrm{b}}^{2}} t^{-\lambda-\rho}+|G|_{\mathcal{C}_{\mathrm{b}}^{1}}\right)\left\|A_{m}^{-\lambda} P_{m} A^{\lambda} \phi\right\|\left\|A_{m}^{-\rho} P_{m} A^{\rho} \psi\right\| \\
& \lesssim t^{-\lambda-\rho}\|\phi\|\|\psi\| .
\end{aligned}
$$

Thus (4.2) is valid.

## 5. Proof of Theorem 1.1

The error will split into several terms, some of which are common to Assumptions $\mathbf{A}$ and $\mathbf{B}$ and some are not. We will first present the proof under Assumption A. When doing so we write it as if the noise were multiplicative, i.e., with the operator $g$ included. This will ease the presentation of the white noise case $\mathbf{B}$.
5.1. The case of colored noise. For an $\mathcal{F}_{T}$-measurable, $H_{m}$-valued random variable $\xi$, the law of iterated expectation and Proposition 1.12 in [3] yields

$$
\begin{equation*}
\mathbf{E}[G(\xi)]=\mathbf{E}\left[\mathbf{E}\left[G(\xi) \mid \mathcal{F}_{T}\right]\right]=\mathbf{E}\left[\mathbf{E}\left[G\left(X^{\xi}(0)\right) \mid \mathcal{F}_{T}\right]\right]=\mathbf{E}\left[u_{m}(0, \xi)\right] \tag{5.1}
\end{equation*}
$$

Thus, the weak error splits into four terms:

$$
\begin{aligned}
& \mathbf{E}\left[G(X(T))-G\left(X_{h}(T)\right)\right] \\
&= \mathbf{E}\left[G(X(T))-G\left(X_{m}(T)\right)\right]+\mathbf{E}\left[G\left(X_{m}(t)\right)-G\left(P_{m} X_{h}(T)\right)\right]+\mathbf{E}\left[G\left(P_{m} X_{h}(T)\right)-G\left(X_{h}(T)\right)\right] \\
&= \mathbf{E}\left[G(X(T))-G\left(X_{m}(T)\right)\right]+u_{m}\left(T, X_{0}\right)-u_{m}\left(T, X_{h}(0)\right) \\
&+\mathbf{E}\left[u_{m}\left(T, X_{h}(0)\right)-u_{m}\left(0, X_{h}(T)\right)\right]+\mathbf{E}\left[G\left(P_{m} X_{h}(T)\right)-G\left(X_{h}(T)\right)\right] \\
&= e_{1}^{m}(T)+e_{2}^{m}(T)+e_{3}^{m}(T)+e_{4}^{m}(T) .
\end{aligned}
$$

Our intention is to let $m \rightarrow \infty$ and see that the remaining terms is of the right order. The first one is easy to treat since when we let $m \rightarrow \infty$ the term $e_{1}^{m}(T) \rightarrow 0$ by the low order of weak convergence implied by the strong convergence. The second term $e_{2}^{m}(T)$ is still easy but needs computations. These holds under both of our assumptions. Using the Cauchy-Schwarz inequality, estimate (4.1) with $0 \leq \lambda=\beta-\epsilon<\alpha=1$, and the error estimate (2.16), we obtain for small $\epsilon>0$

$$
\begin{aligned}
e_{2}^{m}(T) & =u_{m}\left(T, X_{0}\right)-u_{m}\left(T, P_{h} X_{0}\right)=\int_{0}^{1} \frac{\mathrm{~d}}{\mathrm{~d} \lambda} u_{m}\left(T, P_{h} X_{0}+\lambda\left(I-P_{h}\right) X_{0}\right) \mathrm{d} \lambda \\
& =\int_{0}^{1}\left\langle A^{\beta-\epsilon} D u_{m}\left(T, P_{h} X_{0}+\lambda\left(I-P_{h}\right) X_{0}\right), P_{m} A^{-\beta+\epsilon}\left(I-P_{h}\right) X_{0}\right\rangle \mathrm{d} \lambda \\
& \leq \int_{0}^{1}\left\|A^{\beta-\epsilon} D u_{m}\left(T, P_{h} X_{0}+\lambda\left(I-P_{h}\right) X_{0}\right)\right\|_{\mathcal{L}}\left\|P_{m}\right\|_{\mathcal{L}}\left\|A^{\epsilon-\beta}\left(I-P_{h}\right)\right\|_{\mathcal{L}}\left\|X_{0}\right\| \mathrm{d} \lambda \\
& \lesssim h^{2 \beta-2 \epsilon} T^{-\beta}|G|_{C_{\mathrm{b}}^{1}}\left\|X_{0}\right\| \lesssim h^{2 \beta-2 \epsilon}, \quad \text { uniformly in } m .
\end{aligned}
$$

Here we used that $A$ and $P_{m}$ commute and that

$$
\left\|A^{\epsilon-\beta}\left(I-P_{h}\right)\right\|_{\mathcal{L}} \leq\left\|\left(A^{\epsilon-\beta}\left(I-P_{h}\right)\right)^{*}\right\|_{\mathcal{L}} \leq\left\|\left(I-P_{h}\right) A^{\epsilon-\beta}\right\|_{\mathcal{L}}
$$

Notice here that we could have got a sharp result with $\epsilon=0$ under Assumption $\mathbf{A}$, in the case $\beta<1$. However, $e_{2}(T)$ does not allow a sharp rate.

We now turn to the third error term $e_{3}^{m}(T)$. For this we need the Markov generator $L_{h}$ of the finite element solution $X_{h}$. It is given by

$$
\left(L_{h} v\right)(x)=\left\langle A_{h} x-P_{h} f(x), D v(x)\right\rangle-\frac{1}{2} \operatorname{Tr}\left(P_{h} g(x) Q g^{*}(x) P_{h} D^{2} v(x)\right), \quad x \in S_{h}
$$

Itô's formula and the Kolmogorov equation gives that

$$
\begin{aligned}
e_{3}^{m}(T) & =-\left.\mathbf{E}\left[u_{m}\left(T-t, X_{h}(t)\right)-u_{m}\left(T-0, X_{h}(0)\right)\right]\right|_{t=T} \\
& =-\mathbf{E}\left[\int_{0}^{T} \dot{u}_{m}\left(T-t, X_{h}(t)\right)+L_{h} u_{m}\left(T-t, X_{h}(t)\right) \mathrm{d} t\right] \\
& =\mathbf{E} \int_{0}^{T}\left(L_{m}-L_{h}\right) u_{m}\left(T-t, X_{h}(t)\right) \mathrm{d} t
\end{aligned}
$$

The error $e_{3}^{m}(T)$ now naturally divides into three terms:

$$
\begin{aligned}
\left|e_{3}^{m}(T)\right| \leq & \left|\mathbf{E} \int_{0}^{T}\left\langle\left(A_{m}-A_{h}\right) X_{h}(t), D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T}\left\langle\left(P_{m}-P_{h}\right) f\left(X_{h}(t)\right), D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
& +\left\lvert\, \frac{1}{2} \mathbf{E} \int_{0}^{T} \operatorname{Tr}\left\{\left[P_{m} g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right) P_{m}-P_{h} g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right) P_{h}\right]\right.\right. \\
& \left.\times D^{2} u_{m}\left(T-t, X_{h}(t)\right)\right\} \mathrm{d} t \mid \\
= & I+J+K
\end{aligned}
$$

The Ritz projection $R_{h}$ can be expressed in the form $R_{h}=A_{h}^{-1} P_{h} A$. Observing this we can write

$$
\begin{aligned}
\left\langle\left(A_{m}-A_{h}\right) X_{h}, D u_{m}\right\rangle & =\left\langle\left(A_{m} P_{h}-P_{h} A_{h}\right) X_{h}, D u_{m}\right\rangle=\left\langle X_{h},\left(P_{h} A_{m}-A_{h} P_{h}\right) D u_{m}\right\rangle \\
& =\left\langle X_{h}, A_{h} P_{h}\left(A_{h}^{-1} P_{h} A_{m}-I\right) D u_{m}\right\rangle=\left\langle X_{h}, A_{h} P_{h}\left(A_{h}^{-1} P_{h} A P_{m}-I\right) D u_{m}\right\rangle \\
& =\left\langle X_{h}, A_{h} P_{h}\left(R_{h}-I\right) P_{m} D u_{m}\right\rangle+\left\langle X_{h}, A_{h} P_{h}\left(P_{m}-I\right) D u_{m}\right\rangle
\end{aligned}
$$

This enables us to rewrite the term $I$ so that we can apply the error estimates (2.15) and (2.19) for $R_{h}$ and $P_{m}$ respectively. We substitute for $X_{h}$ the mild equation (1.2) and treat the terms separately and estimate

$$
\begin{aligned}
I \leq & \left|\mathbf{E} \int_{0}^{T}\left\langle E_{h}(t) P_{h} X_{0}, A_{h} P_{h}\left(R_{h}-I\right) P_{m} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T}\left\langle\int_{0}^{t} E_{h}(t-s) P_{h} f\left(X_{h}(s)\right) \mathrm{d} s, A_{h} P_{h}\left(R_{h}-I\right) P_{m} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T}\left\langle\int_{0}^{t} E_{h}(t-s) P_{h} g\left(X_{h}(s)\right) \mathrm{d} W(s), A_{h} P_{h}\left(R_{h}-I\right) P_{m} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T}\left\langle A_{h} X_{h},\left(P_{m}-I\right) D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
= & I_{1}^{h}+I_{2}^{h}+I_{3}^{h}+I^{m} .
\end{aligned}
$$

We now estimate $I_{1}^{h}$. Let $\epsilon>0$ be small. Using (2.15), (2.13), (2.18), and (4.1) yields

$$
\begin{aligned}
I_{1}^{h}= & \left|\mathbf{E} \int_{0}^{T}\left\langle A_{h}^{1-\epsilon} E_{h}(t) P_{h} X_{0}, A_{h}^{\epsilon} P_{h}\left(R_{h}-I\right) A^{-\beta+\epsilon} P_{m} A^{\beta-\epsilon} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
\leq & \mathbf{E} \int_{0}^{T}\left\|A_{h}^{1-\epsilon} E_{h}(t) P_{h}\right\|_{\mathcal{L}}\left\|X_{0}\right\|\left\|A_{h}^{\epsilon} P_{h}\left(R_{h}-I\right) A^{-\beta+\epsilon}\right\|_{\mathcal{L}}\left\|P_{m}\right\|_{\mathcal{L}} \\
& \times \sup _{x \in H}\left\|A^{\beta-\epsilon} D u_{m}(T-t, x)\right\| \mathrm{d} t \\
\lesssim & h^{2 \beta-4 \epsilon} \int_{0}^{T} t^{-1+\epsilon}(T-t)^{-\beta+\epsilon} \mathrm{d} t|G|_{\mathcal{C}_{\mathrm{b}}^{1}}\left\|X_{0}\right\| \lesssim h^{2 \beta-4 \epsilon}
\end{aligned}
$$

The term $I_{2}^{h}$ is easily estimated as follows:

$$
\begin{aligned}
I_{2}^{h}= & \mid \mathbf{E} \int_{0}^{T}\left\langle\int_{0}^{t} A_{h}^{1-\epsilon} E_{h}(t-s) P_{h} f\left(X_{h}(s)\right) \mathrm{d} s\right. \\
& \left.A_{h}^{\epsilon} P_{h}\left(R_{h}-I\right) A^{-\beta+\epsilon} P_{m} A^{\beta-\epsilon} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t \mid \\
\leq & \int_{0}^{T} \int_{0}^{t}\left\|A_{h}^{1-\epsilon} E_{h}(t-s) P_{h}\right\|_{\mathcal{L}}\left(\mathbf{E}\left\|f\left(X_{h}(s)\right)\right\|^{2}\right)^{\frac{1}{2}} \\
& \times\left\|A_{h}^{\epsilon} P_{h}\left(R_{h}-I\right) A^{-\beta+\epsilon}\right\|_{\mathcal{L}}\left\|P_{m}\right\|_{\mathcal{L}} \sup _{x \in H}\left\|A^{\beta-\epsilon} D u_{m}(T-t, x)\right\| \mathrm{d} s \mathrm{~d} t .
\end{aligned}
$$

Using (2.15), (2.13), (2.18) and (3.2) yields

$$
I_{2}^{h} \lesssim h^{2 \beta-4 \epsilon} \int_{0}^{T} \int_{0}^{t}(T-t)^{-\beta+\epsilon}(t-s)^{-1+\epsilon} \mathrm{d} s \mathrm{~d} t \lesssim h^{2 \beta-4 \epsilon}
$$

For $I_{3}$ we use the Malliavin integration by parts formula from Lemma 2.2 together with the chain rule (2.25) to obtain the error representation

$$
\begin{aligned}
I_{3}^{h}= & \left|\mathbf{E} \int_{0}^{T}\left\langle\int_{0}^{t} E_{h}(t-s) P_{h} g\left(X_{h}(s)\right) \mathrm{d} W(s), A_{h} P_{h}\left(R_{h}-I\right) P_{m} D u_{m}\left(T-t, X_{h}(t)\right)\right\rangle \mathrm{d} t\right| \\
= & \mid \mathbf{E} \int_{0}^{T} \int_{0}^{t}\left\langle E_{h}(t-s) P_{h} g\left(X_{h}(s)\right),\right. \\
& \left.\quad A_{h} P_{h}\left(R_{h}-I\right) P_{m} D^{2} u_{m}\left(T-t, X_{h}(t)\right) P_{m} D_{s} X_{h}(t)\right\rangle_{\mathcal{L}_{2}^{0}} \mathrm{~d} s \mathrm{~d} t \mid .
\end{aligned}
$$

Here we treat Assumptions $\mathbf{A}$ and $\mathbf{B}$ separately and start with $\mathbf{A} ; \mathbf{B}$ is postponed to the next subsection. Distributing powers of $A$ and $A_{h}$ carefully and setting $g(x)=I$, we write

$$
\begin{gathered}
\left\langle E_{h} P_{h}, A_{h} P_{h}\left(R_{h}-I\right) P_{m} D^{2} u_{m} P_{m} D_{s} X_{h}\right\rangle_{\mathcal{L}_{2}^{0}}=\left\langle A_{h}^{\frac{1+\beta}{2}-\epsilon} E_{h} A_{h}^{\frac{1-\beta}{2}} A_{h}^{\frac{\beta-1}{2}} P_{h}\right. \\
\left.A_{h}^{\frac{1-\beta}{2}+\epsilon} P_{h}\left(R_{h}-I\right) A^{-\frac{1+\beta}{2}+\epsilon} P_{m} A^{\frac{1+\beta}{2}-\epsilon} D^{2} u_{m} A^{\frac{1-\beta}{2}} P_{m} A^{\frac{\beta-1}{2}} D_{s} X_{h}\right\rangle_{\mathcal{L}_{2}^{0}}
\end{gathered}
$$

Using the Cauchy-Schwarz inequality for $\mathcal{L}_{2}^{0}$ and (2.7) yields

$$
\begin{aligned}
I_{3}^{h} \leq \mathbf{E} & \int_{0}^{T} \int_{0}^{t}\left\|A_{h}^{1-\epsilon} E_{h}(t-s) P_{h}\right\|_{\mathcal{L}}\left\|A_{h}^{\frac{\beta-1}{2}} P_{h}\right\|_{\mathcal{L}_{2}^{0}}\left\|A_{h}^{\frac{1-\beta}{2}+\epsilon} P_{h}\left(R_{h}-I\right) A^{-\frac{1+\beta}{2}+\epsilon}\right\|_{\mathcal{L}}\left\|P_{m}\right\|_{\mathcal{L}} \\
& \times \sup _{x \in H}\left\|A^{\frac{1+\beta}{2}-\epsilon} D^{2} u_{m}(T-t, x) A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}}\left\|A^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}} \mathrm{~d} s \mathrm{~d} t .
\end{aligned}
$$

We use (2.13) to get $\left\|A_{h}^{\frac{\beta-1}{2}} P_{h}\right\|_{\mathcal{L}_{2}^{0}} \lesssim\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}$. The norm equivalence (2.14) and the fact that $D_{s}^{u} X_{h}(t) \in S_{h}, \mathbf{P}$-a.s., for every $u \in U_{0}$ yields

$$
\left\|A^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}} \lesssim\left\|A_{h}^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}} .
$$

The analyticity of the semigroup (2.18), the error estimate (2.15) together with (2.13), the gradient estimate (4.2), Tonelli's theorem and the Cauchy-Schwarz inequality now imply that

$$
I_{3}^{h} \lesssim h^{2 \beta-4 \epsilon}|G|_{C_{5}^{2}}\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}} \int_{0}^{T} \int_{0}^{t}\left(\mathbf{E}\left\|A_{h}^{\frac{\beta-1}{2}} D_{s} X_{h}(t)\right\|_{\mathcal{L}_{2}^{0}}^{2}\right)^{\frac{1}{2}}(T-t)^{-1+\epsilon}(t-s)^{-1+\epsilon} \mathrm{d} s \mathrm{~d} t .
$$

Applying Lemma 3.1 we finally get

$$
I_{3}^{h} \lesssim h^{2 \beta-4 \epsilon} \int_{0}^{T} \int_{0}^{t}(T-t)^{-1+\epsilon}(t-s)^{-1+\epsilon} \mathrm{d} s \mathrm{~d} t \lesssim h^{2 \beta-4 \epsilon}
$$

Using (2.19), (4.1), the Cauchy-Schwarz inequality and (2.27) yields

$$
\begin{aligned}
I^{m} & \leq \mathbf{E} \int_{0}^{T}\left\|A_{h} X_{h}(t)\right\|_{\mathcal{L}}\left\|\left(P_{m}-I\right) A^{-\frac{1}{2}+\epsilon}\right\| \sup _{x \in H}\left\|A^{\frac{1}{2}-\epsilon} D u_{m}(T-t, x)\right\|_{\mathcal{L}} \mathrm{d} t \\
& \lesssim \lambda_{m}^{-\frac{1}{2}+\epsilon}\left\|A_{h} P_{h}\right\|_{\mathcal{L}}\left(\sup _{t \in[0, T]} \mathbf{E}\left\|X_{h}(t)\right\|^{2}\right)^{\frac{1}{2}} \int_{0}^{T}(T-t)^{-\frac{1}{2}+\epsilon} \mathrm{d} t .
\end{aligned}
$$

Letting $m \rightarrow \infty$ for fixed $h$ yields $I^{m} \rightarrow 0$ and $\lim _{m \rightarrow \infty} I \lesssim h^{2 \beta-4 \epsilon}$.
The term $J$ is considered next. Writing $P_{m}-P_{h}=\left(P_{m}-I\right)+\left(I-P_{h}\right)$ we get the natural decomposition $J \leq J^{m}+J^{h}$. Using the Cauchy-Schwarz inequality, the error estimate (2.16), (4.1), and (3.2) yields for $i \in\{h, m\}$

$$
\begin{aligned}
J^{i} & =\left|\mathbf{E} \int_{0}^{T}\left\langle\left(I-P_{i}\right) D u_{m}\left(T-t, P_{m} X_{h}(t)\right), f\left(X_{h}(t)\right)\right\rangle \mathrm{d} s\right| \\
& \leq \int_{0}^{T}\left\|\left(I-P_{i}\right) A^{-\beta+\epsilon}\right\|_{\mathcal{L}} \sup _{x \in H^{m}}\left\|A^{\beta-\epsilon} D u_{m}(T-t, x)\right\|\left(\mathbf{E}\left\|f\left(X_{h}(t)\right)\right\|^{2}\right)^{\frac{1}{2}} \mathrm{~d} t \\
& \lesssim\left\|\left(I-P_{i}\right) A^{-\beta+\epsilon}\right\|_{\mathcal{L}}|G|_{\mathcal{C}_{\mathrm{b}}^{1}} \int_{0}^{T}(T-t)^{-\beta+\epsilon} \mathrm{d} t .
\end{aligned}
$$

We have

$$
J^{h} \lesssim h^{2 \beta-2 \epsilon}, \quad \text { and } \quad J^{m} \lesssim \lambda_{m}^{-\beta+\epsilon}
$$

For $K$ we write

$$
\begin{aligned}
& P_{m} g Q g^{*} P_{m}-P_{h} g Q g^{*} P_{h} \\
& \quad=P_{h} g Q g^{*}\left(I-P_{h}\right)+\left(I-P_{h}\right) g Q g^{*} P_{m}+\left(P_{m}+P_{h}\right) g Q g^{*}\left(P_{m}-I\right)
\end{aligned}
$$

and hence we get the following decomposition:

$$
\begin{aligned}
2 K= & \mid \mathbf{E} \int_{0}^{T} \operatorname{Tr}\left(\left[P_{m} g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right) P_{m}-P_{h} g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right) P_{h}\right]\right. \\
& \left.\times D^{2} u_{m}\left(T-t, P_{m} X_{h}(t)\right)\right) \mathrm{d} t \mid \\
\leq & \left|\mathbf{E} \int_{0}^{T} \operatorname{Tr}\left(P_{h} g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right)\left(I-P_{h}\right) D^{2} u_{m}\left(T-t, P_{m} X_{h}(t)\right)\right) \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T} \operatorname{Tr}\left(\left(I-P_{h}\right) g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right) P_{m} D^{2} u_{m}\left(T-t, P_{m} X_{h}(t)\right)\right) \mathrm{d} t\right| \\
& +\left|\mathbf{E} \int_{0}^{T} \operatorname{Tr}\left(\left(P_{m}+P_{h}\right) g\left(X_{h}(t)\right) Q g^{*}\left(X_{h}(t)\right)\left(P_{m}-I\right) D^{2} u_{m}\left(T-t, P_{m} X_{h}(t)\right)\right) \mathrm{d} t\right| \\
= & K_{1}^{h}+K_{2}^{h}+K^{m} .
\end{aligned}
$$

Assumption $\mathbf{A}$ is treated first; $\mathbf{B}$ is postponed. By (2.3), (2.2) and (2.7), we have

$$
\begin{aligned}
& \operatorname{Tr}\left(P_{h} Q\left(I-P_{h}\right) D^{2} u_{m}\right) \\
& \quad=\operatorname{Tr}\left(P_{h} Q\left(I-P_{h}\right) D^{2} u_{m} A^{\frac{1-\beta}{2}} A^{\frac{\beta-1}{2}}\right)=\operatorname{Tr}\left(A^{\frac{\beta-1}{2}} P_{h} Q\left(I-P_{h}\right) D^{2} u_{m} A^{\frac{1-\beta}{2}}\right) \\
& \quad=\operatorname{Tr}\left(A^{\frac{\beta-1}{2}} P_{h} A^{\frac{1-\beta}{2}} A^{\frac{\beta-1}{2}} Q A^{\frac{\beta-1}{2}} A^{\frac{1-\beta}{2}}\left(I-P_{h}\right) A^{-\frac{1+\beta}{2}+\epsilon} A^{\frac{1+\beta}{2}-\epsilon} D^{2} u_{m} A^{\frac{1-\beta}{2}}\right) \\
& \quad \leq\left\|A^{\frac{\beta-1}{2}} P_{h} A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}}\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2}\left\|A^{\frac{1-\beta}{2}}\left(I-P_{h}\right) A^{-\frac{1+\beta}{2}+\epsilon}\right\|_{\mathcal{L}}\left\|A^{\frac{1+\beta}{2}-\epsilon} D^{2} u_{m} A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}},
\end{aligned}
$$

where we used the fact that

$$
\left\|A^{\frac{\beta-1}{2}} Q A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{1}}=\operatorname{Tr}\left(\left(A^{\frac{\beta-1}{2}} Q^{\frac{1}{2}}\right)\left(A^{\frac{\beta-1}{2}} Q^{\frac{1}{2}}\right)^{*}\right)=\left\|A^{\frac{\beta-1}{2}} Q^{\frac{1}{2}}\right\|_{\mathcal{L}_{2}}^{2}=\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2}
$$

By (2.14) and (2.13) $\left\|A^{\frac{\beta-1}{2}} P_{h} A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}} \lesssim\left\|A_{h}^{\frac{\beta-1}{2}} P_{h} A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}} \lesssim\left\|A^{\frac{\beta-1}{2}} A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}}=1 . \quad$ Using (2.16), (3.2) and (4.2) gives us

$$
K_{1}^{h} \lesssim h^{2 \beta-2 \epsilon}\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2}|G|_{\mathcal{C}_{\mathrm{b}}^{2}} \int_{0}^{T}(T-t)^{-1+\epsilon} \mathrm{d} t \lesssim h^{2 \beta-2 \epsilon}
$$

For $K_{2}^{h}$ we compute similarly

$$
\begin{aligned}
\operatorname{Tr}\left(\left(I-P_{h}\right) Q P_{m} D^{2} u\right) & =\operatorname{Tr}\left(A^{-\frac{1+\beta}{2}+\epsilon}\left(I-P_{h}\right) A^{\frac{1-\beta}{2}} A^{\frac{\beta-1}{2}} Q A^{\frac{\beta-1}{2}} A^{\frac{1-\beta}{2}} D^{2} u_{m} A^{\frac{1+\beta}{2}-\epsilon}\right) \\
& \leq\left\|A^{-\frac{1+\beta}{2}+\epsilon}\left(I-P_{h}\right) A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}}\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2}\left\|A^{\frac{1-\beta}{2}} D^{2} u_{m} A^{\frac{1+\beta}{2}-\epsilon}\right\|_{\mathcal{L}}
\end{aligned}
$$

where

$$
\left\|A^{-\frac{1+\beta}{2}+\epsilon}\left(I-P_{h}\right) A^{\frac{1-\beta}{2}}\right\|_{\mathcal{L}} \leq\left\|\left(A^{-\frac{1+\beta}{2}+\epsilon}\left(I-P_{h}\right) A^{\frac{1-\beta}{2}}\right)^{*}\right\|_{\mathcal{L}}=\left\|A^{\frac{1-\beta}{2}}\left(I-P_{h}\right) A^{-\frac{1+\beta}{2}+\epsilon}\right\|_{\mathcal{L}}
$$

so that (2.16) applies. Hence,

$$
K_{2}^{h} \lesssim h^{2 \beta-2 \epsilon}\left\|A^{\frac{\beta-1}{2}}\right\|_{\mathcal{L}_{2}^{0}}^{2}|G|_{C_{\mathrm{b}}^{2}} \int_{0}^{T}(T-t)^{-1+\epsilon} \mathrm{d} t \lesssim h^{2 \beta-2 \epsilon}
$$

Term $K^{m}$ is treated analogously as $K_{1}^{h}$. We have $K^{m} \lesssim \lambda_{m}^{-\beta+\epsilon}$.
Finally, by the Lipschitz continuity of $G$, the Dominated Convergence Theorem and the strong convergence of $P_{m} \rightarrow I$ we get

$$
e_{4}^{m}(T) \leq|G|_{\mathcal{C}_{\mathrm{b}}^{1}} \mathbf{E}\left\|\left(P_{m}-I\right) X_{h}(t)\right\| \rightarrow 0, \quad \text { as } m \rightarrow \infty
$$

We conclude that $e(T)=O\left(h^{2 \gamma}\right)$ for any $\gamma<\beta$, which completes the proof in case $\mathbf{A}$.
5.2. The case of white noise. Now consider the case of Assumption B. All estimates above, except those for $I_{3}^{h}, K_{1}^{h}, K_{2}^{h}$ and $K^{m}$, hold under $\mathbf{B}$ by setting $Q=I$ and $\beta=\frac{1}{2}$ and recalling that $U_{0}=H$ and $\mathcal{L}_{2}^{0}=\mathcal{L}_{2}$. We now complete the proof with the remaining estimates.

Using Hölder's inequality (2.8) yields

$$
\begin{aligned}
I_{3}^{h}=\mid \mathbf{E} & \int_{0}^{T} \int_{0}^{t}\left\langle E_{h}(t-s) P_{h} g\left(X_{h}(s)\right), A_{h} P_{h}\left(R_{h}-I\right)\right. \\
& \left.\times P_{m} D^{2} u_{m}\left(T-t, P_{m} X_{h}(t)\right) P_{m} D_{s} X_{h}(t)\right\rangle_{\mathcal{L}_{2}} \mathrm{~d} s \mathrm{~d} t \mid \\
\leq & \mathbf{E} \int_{0}^{T} \int_{0}^{t}\left\|A_{h}^{1-3 \epsilon} E_{h}(t-s) P_{h}\right\|_{\mathcal{L}}\left\|g\left(X_{h}(s)\right)\right\|_{\mathcal{L}}\left\|A_{h}^{3 \epsilon} P_{h}\left(R_{h}-I\right) A^{-\frac{1}{2}+\epsilon}\right\|_{\mathcal{L}} \\
& \times \sup _{x \in H^{m}}\left\|A^{\frac{1}{2}-\epsilon} D^{2} u_{m}(T-t, x) A^{\frac{1}{2}-\epsilon}\right\|_{\mathcal{L}}\left\|A^{-\frac{1}{2}+\epsilon} A_{h}^{-2 \epsilon} P_{h}\right\|_{\mathcal{L}_{1}}\left\|A_{h}^{2 \epsilon} D_{s} X_{h}(t)\right\|_{\mathcal{L}} \mathrm{d} s \mathrm{~d} t
\end{aligned}
$$

First, using (2.13) and (2.9), we have

$$
\left\|A^{-\frac{1}{2}+\epsilon} A_{h}^{-2 \epsilon} P_{h}\right\|_{\mathcal{L}_{1}}=\left\|A_{h}^{-2 \epsilon} P_{h} A^{-\frac{1}{2}+\epsilon}\right\|_{\mathcal{L}_{1}} \lesssim\left\|A_{h}^{-2 \epsilon} P_{h} A^{2 \epsilon}\right\|_{\mathcal{L}}\left\|A^{-\frac{1}{2}-\epsilon}\right\|_{\mathcal{L}_{1}} \lesssim\left\|A^{-\frac{1}{2}-\epsilon}\right\|_{\mathcal{L}_{1}}
$$

Now we apply (2.15), (2.18), (4.2) with $\rho=\lambda=\frac{1}{2}-\epsilon<\alpha=\frac{1}{2}$, to get

$$
I_{3}^{h} \lesssim h^{1-8 \epsilon}|G|_{\mathcal{C}_{\mathrm{b}}^{2}} \int_{0}^{T} \int_{0}^{t}\left(\mathbf{E}\left\|g\left(X_{h}(s)\right)\right\|_{\mathcal{L}}^{2}\right)^{\frac{1}{2}}\left(\mathbf{E}\left\|A_{h}^{2 \epsilon} D_{s} X_{h}(t)\right\|_{\mathcal{L}}^{2}\right)^{\frac{1}{2}}(T-t)^{-1+2 \epsilon}(t-s)^{-1+3 \epsilon} \mathrm{~d} s \mathrm{~d} t
$$

Finally using Lemma 3.2 together with (3.2) finishes the estimate of $I_{3}^{h}$. Indeed,

$$
I_{3}^{h} \lesssim h^{1-8 \epsilon}|G|_{C_{\mathrm{b}}^{2}} \int_{0}^{T}(T-t)^{-1+2 \epsilon}(t-s)^{-1+\epsilon} \mathrm{d} s \mathrm{~d} t \lesssim h^{1-8 \epsilon}
$$

For $K_{1}$ we use Hölder's inequality (2.6), (2.16), and Lemma 4.2 to get

$$
\begin{aligned}
2 K_{1} \leq & \int_{0}^{T} \mathbf{E}\left\|A^{-\frac{1-\epsilon}{2}} P_{h} g\left(X_{h}(t)\right) g^{*}\left(X_{h}(t)\right) A^{-\epsilon}\right\|_{\mathcal{L}_{1}}\left\|A^{\epsilon}\left(I-P_{h}\right) A^{-\frac{1-\epsilon}{2}}\right\|_{\mathcal{L}} \\
& \times \sup _{x \in H}\left\|A^{\frac{1-\epsilon}{2}} D^{2} u_{m}(T-t, x) A^{\frac{1-\epsilon}{2}}\right\|_{\mathcal{L}} \mathrm{d} t \\
\lesssim & h^{1-3 \epsilon} \sup _{t \in[0, T]} \mathbf{E}\left\|A^{-\frac{1-\epsilon}{2}} P_{h} g\left(X_{h}(t)\right) g^{*}\left(X_{h}(t)\right)\right\|_{\mathcal{L}_{2 /(2-3 \epsilon)}}\left\|A^{-\epsilon}\right\|_{\mathcal{L}_{2 / 3 \epsilon}}|G|_{\mathcal{C}_{\mathrm{b}}^{2}} \int_{0}^{T}(T-t)^{-1+\epsilon} \mathrm{d} t \\
\lesssim & h^{1-3 \epsilon} \sup _{t \in[0, T]} \mathbf{E}\left\|g\left(X_{h}(t)\right)\right\|_{\mathcal{L}}^{2}\left\|A^{-\frac{1-\epsilon}{2}}\right\|_{\mathcal{L}_{2 /(2-3 \epsilon)}}\left\|A^{-\epsilon}\right\|_{\mathcal{L}_{2 / 3 \epsilon}}|G|_{\mathcal{C}_{\mathrm{b}}^{2}} .
\end{aligned}
$$

We compute and use (2.9) to conclude

$$
\begin{aligned}
\left\|A^{-\epsilon}\right\|_{\mathcal{L}_{2 / 3 \epsilon}}^{3 \epsilon / 2} & =\sum_{i \in \mathbf{N}}\left(\lambda_{i}^{-\epsilon}\right)^{\frac{2}{3 \epsilon}}=\sum_{i \in \mathbf{N}} \lambda_{i}^{-\frac{2}{3}}=\operatorname{Tr}\left(A^{-\frac{2}{3}}\right)<\infty \\
\left\|A^{-\frac{1-\epsilon}{2}}\right\|_{\mathcal{L}_{2 /(2-3 \epsilon)}}^{(2-3 \epsilon) / 2} & =\sum_{i \in \mathbf{N}}\left(\lambda_{i}^{-\frac{1-\epsilon}{2}}\right)^{\frac{2}{2-3 \epsilon}}=\sum_{i \in \mathbf{N}} \lambda_{i}^{-\frac{1-\epsilon}{2-3 \epsilon}}=\operatorname{Tr}\left(A^{-\frac{1}{2}\left(\frac{2-2 \epsilon}{2-3 \epsilon}\right)}\right)<\infty
\end{aligned}
$$

The terms $K_{2}^{h}$ and $K^{m}$ admits the same treatment, so that

$$
K_{2}^{h} \lesssim h^{1-3 \epsilon}, \quad \text { and } \quad K^{m} \lesssim \lambda_{m}^{-\frac{1}{2}+\frac{3 \epsilon}{2}}
$$

We have $e(T)=\mathcal{O}\left(h^{2 \gamma}\right)$ for any $\gamma<\frac{1}{2}$.

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