

Strong uniqueness for stochastic evolution equations with unbounded measurable drift term

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Abstract We consider stochastic evolution equations in Hilbert spaces with merely measurable and locally bounded drift term B and cylindrical Wiener noise. We prove pathwise (hence strong) uniqueness in the class of global solutions. This paper extends our previous paper [5] which generalized Veretennikov’s fundamental result to infinite dimensions assuming boundedness of the drift term. As in [5] pathwise uniqueness holds for a large class, but not for every initial condition. We also include an application of our result to prove existence of strong solutions when the drift B is only measurable, locally bounded and grows more than linearly.

Key words: Pathwise uniqueness, stochastic PDEs, locally bounded measurable drift term, strong mild solutions

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1 Introduction

We consider the following abstract stochastic differential equation in a separable Hilbert space H

$$dX_t = (AX_t + B(X_t))dt + dW_t, \quad t \geq 0, \quad X_0 = x \in H, \quad (1)$$

where $A : D(A) \subset H \rightarrow H$ is self-adjoint, negative definite and such that $(-A)^{-1+\delta}$, for some $\delta \in (0, 1)$, is of trace class, $B : H \rightarrow H$ and $W = (W_t)$ is a cylindrical Wiener process. In [5] under the assumption that B is Borel measurable and (globally) bounded we prove pathwise uniqueness of solutions to (1). A natural generalization is to extend it to the case where we only assume that B is *Borel measurable and locally bounded* (i.e., bounded on balls):

$$B \in B_{b,loc}(H, H). \quad (2)$$

In this paper we prove that assuming (2) pathwise uniqueness holds for μ -a.e. initial condition x in the class of global mild solutions to (1). Here μ denotes the Gaussian measure which is invariant for the Ornstein-Uhlenbeck process $Z = (Z_t)$ which solves (1) when $B = 0$ (see Section 1.1 for more details).

In other words if for some initial condition $x \in H$, μ -a.e., there exists a solution for (1) on some filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ with a cylindrical (\mathcal{F}_t) -Wiener process W then our main result shows that this solution is pathwise unique. This is in particular the case when

$$B \text{ is measurable and at most of linear growth} \quad (3)$$

(i.e., B is measurable and there exists $a, b \geq 0$ such that $|B(x)| \leq a + b|x|$, $x \in H$), because then existence of weak mild solutions is well-known (see Chapter 10 in [6], [11, 14] and also Appendix in [5]). Moreover, under condition (3), the unique law of any mild solution X^x is equivalent to the law of the Ornstein-Uhlenbeck process starting at x (corresponding to $B = 0$).

By our main result, using a generalization of the Yamada-Watanabe theorem (see [18, 21]), one deduces that under (3) equation (1) has a unique strong mild solution, for μ -a.e. $x \in H$, generalizing A. Veretennikov's seminal result [26] in the case $H = \mathbb{R}^d$ (see also [9, 10, 15, 16, 17, 25, 27, 28]) to infinite dimensions.

In Section 4 we generalize this by relaxing assumption (3). We prove existence of strong mild solutions, starting from μ -a.e. initial condition $x \in H$, when $B \in B_{b,loc}(H, H)$ and moreover there exist $C > 0$, $p > 0$, such that

$$\langle B(y+z), y \rangle \leq C(|y|^2 + e^{p|z|} + 1) \quad (4)$$

for all $y, z \in H$ (see also Remark 17). Finally in Section 4.3 we show a possible extension of our result by considering local mild solutions.

In order to prove *pathwise uniqueness* for (1) we will consider bounded truncated drifts like

$$B_N = B 1_{B(0,N)}, \quad N \geq 1, \quad (5)$$

where $B(0, N)$ is the open ball of center 0 and of radius N and $1_{B(0, N)}$ is the indicator function of $B(0, N)$, by performing a suitable stopping time argument. This argument is not straightforward since it must be also used in combination with the Ito-Tanaka trick from [5] (see also [8, 12, 13, 23, 26]).

In addition in this paper we also simplify some arguments used in [5] in the case of $B \in B_b(H, H)$ (see, in particular, Lemma 8). Before stating our main result precisely, let us recall the following definition (cf. [21] and [18]).

Definition 1. *Let $x \in H$.*

(a) *We call weak mild solution to (1) a tuple $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}, W, X)$, where $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ is a filtered probability space on which it is defined a cylindrical (\mathcal{F}_t) -Wiener process W and a continuous (\mathcal{F}_t) -adapted H -valued process $X = (X_t) = (X_t)_{t \geq 0}$ such that, \mathbb{P} -a.s.,*

$$X_t = e^{tA}x + \int_0^t e^{(t-s)A}B(X_s)ds + \int_0^t e^{(t-s)A}dW_s, \quad t \geq 0. \quad (6)$$

(b) *A weak mild solution X which is $(\bar{\mathcal{F}}_t^W)$ -adapted (here $(\bar{\mathcal{F}}_t^W)$ denotes the completed natural filtration of the cylindrical process W) is called strong mild solution.*

We will often use stopping times

$$\tau_N^X = \inf\{t \geq 0 : X_t \notin B(0, N)\} \quad (7)$$

($\tau_N^X = +\infty$ if the set is empty), $N \geq 1$.

In our main result we consider two mild solutions X and Y , having the same initial condition $x \in H$, and solving the same equation (1) but with possibly different drift terms, respectively B and $B' \in B_{b,loc}(H, H)$, i.e.,

$$dX_t = (AX_t + B(X_t))dt + dW_t, \quad X_0 = x, \quad (8)$$

$$dY_t = (AY_t + B'(Y_t))dt + dW_t, \quad Y_0 = x. \quad (9)$$

Theorem 1. *Assume Hypothesis 1 (see Section 1.1) and let μ be the centered Gaussian measure on H with covariance $Q = -\frac{1}{2}A^{-1}$.*

Then for μ -a.e. $x \in H$, if X and Y are two weak mild solutions, respectively of (8) and (9), defined on the same filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ with the same cylindrical Wiener process W , and if, for some $N \geq 1$,

$$B(x) = B'(x), \quad x \in B(0, N), \quad (10)$$

then, \mathbb{P} -a.s.,

$$X_{t \wedge \tau_N^X \wedge \tau_N^Y} = Y_{t \wedge \tau_N^X \wedge \tau_N^Y}, \quad t \geq 0, \quad (11)$$

and so $\tau_N^X = \tau_N^Y$, \mathbb{P} -a.s..

Above we restrict to W which are cylindrical with respect to the eigenbasis of A (see Section 1.1 for details). Clearly if $B = B'$ the result implies that, \mathbb{P} -a.s.,

$$X_t = Y_t, \quad t \geq 0. \quad (12)$$

Indeed using that $\tau_N^X \uparrow +\infty$ and $\tau_N^Y \uparrow +\infty$ as $N \rightarrow \infty$ (because X and Y are both global solutions) we deduce easily (12) from (11).

The proof of Theorem 1, performed in Section 3, uses a truncation argument and regularity results for elliptic equations in Hilbert spaces involving truncated drift terms B_N (cf. (5)). Such regularity results are given in Section 2, where we also establish an Itô type formula involving $u(X_t)$ with u in some Sobolev space associated to μ (see Theorem 10). In comparison with [5] to prove such an Itô type formula we use a new analytic lemma (see Lemma 8).

There are several other quite essential differences in comparison with [5] in our proof. We refer to Remarks 9 and 12 for details.

1.1 Assumptions and preliminaries

As in [5] we are given a real separable Hilbert space H and denote its norm and inner product by $|\cdot|$ and $\langle \cdot, \cdot \rangle$ respectively. We follow [4, 6, 7, 22], and assume

Hypothesis 1 $A : D(A) \subset H \rightarrow H$ is a negative definite self-adjoint operator and $(-A)^{-1+\delta}$, for some $\delta \in (0, 1)$, is of trace class.

Remark 2. Our uniqueness result continues to hold under the following more general assumption: $A : D(A) \subset H \rightarrow H$ is self-adjoint and there exists $\omega \in \mathbb{R}$ such that $(A - \omega)$ is negative definite and $(\omega - A)^{-1+\delta}$, for some $\delta \in (0, 1)$, is of trace class. Indeed if we write equation (1) in the form

$$dX_t = (AX_t - \omega X_t)dt + (\omega X_t + B(X_t))dt + dW_t, \quad X_0 = x \in H,$$

then the linear operator $(A - \omega I)$ verifies Hypothesis 1 and the drift $\omega I + B$ continues to satisfy (2).

Since A^{-1} is compact, there exists an orthonormal basis (e_k) in H and a sequence of positive numbers (λ_k) such that

$$Ae_k = -\lambda_k e_k, \quad k \in \mathbb{N}. \quad (13)$$

Recall that A generates an analytic semigroup e^{tA} on H such that $e^{tA}e_k = e^{-\lambda_k t}e_k$.

From now on until and including Section 3 we fix $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$, W , X and Y as in the assertion of Theorem 1. As said before we will consider a cylindrical Wiener process W_t with respect to the previous basis (e_k) . The process W_t is formally given by “ $W_t = \sum_{k \geq 1} \beta_k(t)e_k$ ” where $\beta_k(t)$ are independent one dimensional Wiener processes (see [6] for more details).

By R_t we denote the Ornstein-Uhlenbeck semigroup in $B_b(H)$ (the Banach space of Borel and bounded real functions endowed with the essential supremum norm $\|\cdot\|_0$) defined as

$$R_t\varphi(x) = \int_H \varphi(y)N(e^{tA}x, Q_t)(dy), \quad \varphi \in B_b(H), \quad (14)$$

where $N(e^{tA}x, Q_t)$ is the Gaussian measure in H of mean $e^{tA}x$ and covariance operator Q_t given by,

$$Q_t = -\frac{1}{2} A^{-1}(I - e^{2tA}), \quad t \geq 0. \quad (15)$$

We note that R_t has the unique invariant measure $\mu := N(0, Q)$ where $Q = -\frac{1}{2} A^{-1}$. Moreover, since under the previous assumptions, the Ornstein-Uhlenbeck semigroup is strong Feller and irreducible we have by Doob's theorem that, for any $t > 0$, $x \in H$, the measures $N(e^{tA}x, Q_t)$ and μ are equivalent (see [7]). On the other hand, our assumption that $(-A)^{-1+\delta}$ is trace class guarantees that the OU process

$$Z_t = Z(t, x) = e^{tA}x + \int_0^t e^{(t-s)A}dW_s \quad (16)$$

has a continuous H -valued version.

If H and K are separable Hilbert spaces, the Banach space $L^p(H, \mu; K)$, $p \geq 1$, is defined to consist of equivalent classes of measurable functions $f : H \rightarrow K$ such that $\int_H |f|_K^p \mu(dx) < +\infty$ (if $K = \mathbb{R}$ we set $L^p(H, \mu; \mathbb{R}) = L^p(H, \mu)$). We also use the notation $L^p(\mu)$ instead of $L^p(H, \mu, K)$ when no confusion may arise.

The semigroup R_t can be uniquely extended to a strongly continuous semigroup of contractions on $L^p(H, \mu)$, $p \geq 1$, which we still denote by R_t , whereas we denote by L_p (or L when no confusion may arise) its infinitesimal generator, which is defined on smooth functions φ as

$$L\varphi(x) = \frac{1}{2}\text{Tr}(D^2\varphi(x)) + \langle Ax, D\varphi(x) \rangle,$$

where $D\varphi(x)$ and $D^2\varphi(x)$ denote respectively the first and second Fréchet derivatives of φ at $x \in H$. For Banach spaces E and F we denote by $C_b^k(E, F)$, $k \geq 1$, the Banach space of all functions $f : E \rightarrow F$ which are bounded and Fréchet differentiable on E up to order $k \geq 1$ with all derivatives bounded and continuous. We also set $C_b^k(E, \mathbb{R}) = C_b^k(E)$.

According to [7], for any $\varphi \in B_b(H)$ and any $t > 0$ one has $R_t\varphi \in C_b^\infty(H) = \bigcap_{k \geq 1} C_b^k(H)$. Moreover,

$$\langle DR_t\varphi(x), h \rangle = \int_H \langle \Lambda_t h, Q_t^{-\frac{1}{2}} y \rangle \varphi(e^{tA}x + y) N(0, Q_t)(dy), \quad h \in H, \quad (17)$$

where Q_t is as defined in (15),

$$\Lambda_t = Q_t^{-1/2} e^{tA} = \sqrt{2} (-A)^{1/2} e^{tA} (I - e^{2tA})^{-1/2} \quad (18)$$

and $y \mapsto \langle \Lambda_t h, Q_t^{-\frac{1}{2}} y \rangle$ is a centered Gaussian random variable under $\mu_t = N(0, Q_t)$ with variance $|\Lambda_t h|^2$ for any $t > 0$ (cf. Theorem 6.2.2 in [6]). Since

$$\Lambda_t e_k = \sqrt{2} (\lambda_k)^{1/2} e^{-t\lambda_k} (1 - e^{-2t\lambda_k})^{-1/2} e_k,$$

we see that there exists $C'_0 > 0$ such that $\|\Lambda_t\| \leq C'_0 t^{-\frac{1}{2}}$.

In the sequel $\|\cdot\|$ always denotes the *Hilbert-Schmidt norm*; on the other hand $\|\cdot\|_{\mathcal{L}}$ indicates the *operator norm*. By (17) we deduce

$$\sup_{x \in H} |DR_t \varphi(x)| = \|DR_t \varphi\|_0 \leq C_0 t^{-\frac{1}{2}} \|\varphi\|_0, \quad t > 0, \quad (19)$$

which by taking the Laplace transform yields

$$\|D(\lambda - L_2)^{-1} \varphi\|_0 \leq \frac{C_0}{\lambda^{\frac{1}{2}}} \|\varphi\|_0, \quad \lambda > 0.$$

Similarly, we find

$$\|DR_t \varphi\|_{L^2(\mu)} \leq C_0 t^{-\frac{1}{2}} \|\varphi\|_{L^2(\mu)}$$

and

$$\|D(\lambda - L_2)^{-1} \varphi\|_{L^2(\mu)} \leq \frac{C_0}{\lambda^{\frac{1}{2}}} \|\varphi\|_{L^2(\mu)}.$$

Recall that the Sobolev space $W^{2,p}(H, \mu)$, $p \geq 1$, is defined in Section 3 of [3] as the completion of a suitable set of smooth functions endowed with the Sobolev norm (see also Section 9.2 in [6] for the case $p = 2$ and [24]). Under our above assumptions, the following result can be found in Section 10.2.1 of [7].

Theorem 3. *Let $\lambda > 0$, $f \in L^2(H, \mu)$ and let $\varphi \in D(L_2)$ be the solution of the equation*

$$\lambda \varphi - L_2 \varphi = f.$$

Then $\varphi \in W^{2,2}(H, \mu)$, $(-A)^{1/2} D\varphi \in L^2(H, \mu; H)$ and there exists a constant $C(\lambda)$ such that

$$\|\varphi\|_{L^2(\mu)} + \left(\int_H \|D^2 \varphi(x)\|^2 \mu(dx) \right)^{1/2} + \|(-A)^{1/2} D\varphi\|_{L^2(\mu)} \leq C \|f\|_{L^2(\mu)}.$$

The following extension to $L^p(\mu)$, $p > 1$, can be found in Section 3 of [3] (see also [2, 20]).

Theorem 4. *Let $\lambda > 0$, $f \in L^p(H, \mu)$ and let $\varphi \in D(L_p)$ be the solution of the equation*

$$\lambda \varphi - L_p \varphi = f.$$

Then $\varphi \in W^{2,p}(H, \mu)$, $(-A)^{1/2} D\varphi \in L^p(H, \mu; H)$ and there exists a constant $C = C(\lambda, p)$ such that

$$\|\varphi\|_{L^p(\mu)} + \left(\int_H \|D^2 \varphi(x)\|^p \mu(dx) \right)^{1/p} + \|(-A)^{1/2} D\varphi\|_{L^p(\mu)} \leq C \|f\|_{L^p(\mu)}.$$

2 Analytic results and an Itô type formula

2.1 Existence and uniqueness for the Kolmogorov equation when B is bounded

We are here concerned with the equation

$$\lambda u - L_2 u - \langle B, Du \rangle = f, \quad (20)$$

where $\lambda > 0$, $f \in B_b(H)$ and $B \in B_b(H, H)$ (i.e., $B : H \rightarrow H$ is Borel and bounded).

Remark 5. Since the corresponding Dirichlet form

$$\mathcal{E}(u, v) := \int_H \langle Du, Dv \rangle d\mu - \int_H \langle B, Du \rangle v d\mu + \lambda \int_H uv d\mu,$$

$u, v \in W^{1,2}(\mu)$, is weakly sectorial for λ big enough, it follows by Chap. I and Subsection 3e in Chap. II of [19] that (20) has a unique solution in $D(L_2)$. However, we need more regularity for u .

We recall a result from [5].

Proposition 6. *Let $\lambda \geq \lambda_0$, where*

$$\lambda_0 := 4\|B\|_0^2 C_0^2. \quad (21)$$

Then there is a unique solution $u \in D(L_2)$ of (20) given by

$$u = u_\lambda = (\lambda - L_2)^{-1}(I - T_\lambda)^{-1}f,$$

where

$$T_\lambda \varphi := \langle B, D(\lambda - L_2)^{-1} \varphi \rangle. \quad (22)$$

Moreover, $u \in C_b^1(H)$ with

$$\|u\|_0 \leq 2\|f\|_0, \quad \|Du\|_0 \leq \frac{2C_{1,0}}{\lambda^{\frac{1}{2}}} \|f\|_0, \quad (23)$$

and, for any $p \geq 2$, $u \in W^{2,p}(H, \mu)$ and, for some $C = C(\lambda, p, \|B\|_0)$,

$$\int_H \|D^2 u(x)\|^p \mu(dx) \leq C \int_H |f(x)|^p \mu(dx). \quad (24)$$

2.2 Approximations

We are given two sequences $(f_n) \subset B_b(H)$ and $(B_n) \subset B_b(H, H)$ such that

$$\begin{aligned} (i) \quad & f_n(x) \rightarrow f(x), \quad B_n(x) \rightarrow B(x) \quad \mu\text{-a.e.} \\ (ii) \quad & \|f_n\|_0 \leq M, \quad \|B_n\|_0 \leq M. \end{aligned} \tag{25}$$

The following result has been proved in [5].

Proposition 7. *Let $\lambda \geq \lambda_0$, where λ_0 is defined in (21). Then the equation*

$$\lambda u_n - Lu_n - \langle B_n, Du_n \rangle = f_n, \tag{26}$$

has a unique solution $u_n \in C_b^1(H) \cap D(L_2)$ given by

$$u_n = (\lambda - L)^{-1}(I - T_{n,\lambda})^{-1}f_n,$$

where

$$T_{n,\lambda}\varphi := \langle B_n, D(\lambda - L_2)^{-1}\varphi \rangle.$$

Moreover,

$$\|u_n\|_0 \leq 2M, \quad \|Du_n\|_0 \leq \frac{2C_{1,0}}{\lambda^{\frac{1}{2}}}M, \quad n \geq 1. \tag{27}$$

Finally, we have $u_n \rightarrow u$, and $Du_n \rightarrow Du$, in $L^2(\mu)$, where u is the solution to (20).

Next we prove a new result. The idea behind the result is that if (f_n) satisfies (25) then, for any $x \in H$ (not only μ -a.e.), $t > 0$,

$$R_t f_n(x) \rightarrow R_t f(x) \tag{28}$$

as $n \rightarrow \infty$, due to the fact that, for any $x \in H$, the law of the OU process $Z(t, x)$ at time $t > 0$ is absolutely continuous with respect to μ .

Lemma 8. *Consider the situation of Proposition 7. Then we have:*

$$u_n(x) \rightarrow u(x), \quad Du_n(x) \rightarrow Du(x), \tag{29}$$

for any $x \in H$.

Proof. By a standard argument, possibly passing to a subsequence, we may assume that $Du_n(x) \rightarrow Du(x)$, μ -a.e.. It follows that for any x , μ -a.e.,

$$f_n(x) + \langle B_n(x), Du_n(x) \rangle \rightarrow f(x) + \langle B(x), Du(x) \rangle,$$

as $n \rightarrow \infty$. We write, for any $\lambda \geq \lambda_0$, $x \in H$,

$$u_n(x) = \int_0^\infty e^{-\lambda t} R_t (f_n + \langle B_n, Du_n \rangle)(x) dt.$$

By the argument used in (28) we can apply the dominated convergence theorem and obtain that $u_n(x) \rightarrow u(x)$, as $n \rightarrow \infty$, $x \in H$.

Concerning Du_n , using (19), we obtain, for any $x \in H$, $h \in H$,

$$\langle Du_n(x), h \rangle = \int_0^\infty e^{-\lambda t} \langle DR_t(f_n + \langle B_n, Du_n \rangle)(x), h \rangle dt.$$

Setting $g_n = f_n + \langle B_n, Du_n \rangle$, $g = f + \langle B, Du \rangle$, note that

$$\langle DR_t g_n(x), e_k \rangle = \int_H \langle \Lambda_t e_k, Q_t^{-\frac{1}{2}} y \rangle g_n(e^{tA} x + y) N(0, Q_t)(dy), \quad k \geq 1, x \in H, t > 0.$$

It follows that

$$|\langle DR_t(g_n - g)(x), e_k \rangle|^2 \leq |\Lambda_t e_k|^2 \int_H |g(e^{tA} x + y) - g_n(e^{tA} x + y)|^2 N(0, Q_t)(dy),$$

and so

$$|DR_t(g_n - g)(x)|^2 \leq \|\Lambda_t\|^2 \int_H |g(e^{tA} x + y) - g_n(e^{tA} x + y)|^2 N(0, Q_t)(dy).$$

Now we get $|DR_t(g_n - g)(x)|^2 \rightarrow 0$ as $n \rightarrow \infty$, for any $x \in H$, $t > 0$, by the same argument used in (28).

Using again the dominated convergence theorem we obtain easily that $Du_n(x) \rightarrow Du(x)$, as $n \rightarrow \infty$, for any $x \in H$. \square

2.3 Modified mild formulation

Recall the notation

$$B_N = B \mathbf{1}_{B(0, N)}, \quad N \geq 1, \quad (30)$$

where $B(0, N)$ is the open ball of radius N (hence $B_N \in B_b(H, H)$, $N \geq 1$).

For any $i \in \mathbb{N}$ we denote the i^{th} component of B by $B^{(i)}$, i.e.,

$$B^{(i)}(x) := \langle B(x), e_i \rangle, \quad x \in H.$$

Then for $\lambda \geq 4\|B_N\|_0^2 C_0^2$ we consider the solution $u_N^{(i)}$ of the equation

$$\lambda u_N^{(i)} - Lu_N^{(i)} - \langle B_N, Du_N^{(i)} \rangle = B_N^{(i)}, \quad \mu \text{-a.e.} \quad (31)$$

We recall that by Proposition 6, $u_N^{(i)} \in C_b^1(H)$ and, for any $p \geq 2$, $u_N^{(i)} \in W^{2,p}(H, \mu)$.

The next result is a kind of “local version” of Theorem 7 in [5]. In contrast to Theorem 1 the result holds for any initial condition $x \in H$.

Remark 9. Compared with the proof of Theorem 7 in [5], here we will use Lemma 8 which allows to simplify some arguments of [5] and it is also needed to justify the approximation procedure (see in particular (45)). We also mention that differently with respect to [5] in Step 3 of the proof we need to construct a suitable auxiliary process $\hat{X}^N = (\hat{X}_t^N)$ (see (47)) in order to apply the Girsanov theorem and get the assertion.

Theorem 10. *Let $X = (X_t)$ be a weak mild solution of equation (1) defined on some filtered probability space with a cylindrical (\mathcal{F}_t) -Wiener process W . Consider the stopping time*

$$\tau_N^X = \inf\{t \geq 0 : X_t \notin B(0, N)\}$$

Let $u_N^{(i)}$ be the solution of (31) and set $X_t^{(i)} = \langle X_t, e_i \rangle$. For any $t > 0$ we have \mathbb{P} -a.s. on the event $\{t \leq \tau_N^X\}$

$$\begin{aligned} X_t^{(i)} &= e^{-\lambda_i t} (\langle x, e_i \rangle + u_N^{(i)}(x)) - u_N^{(i)}(X_t) \\ &\quad + (\lambda + \lambda_i) \int_0^t e^{-\lambda_i(t-s)} u_N^{(i)}(X_s) ds \\ &\quad + \int_0^t e^{-\lambda_i(t-s)} (d\langle W_s, e_i \rangle + \langle Du_N^{(i)}(X_s), dW_s \rangle). \end{aligned} \quad (32)$$

Proof. We fix $t > 0$, $N \geq 1$ and $i \geq 1$.

Step 1 Approximation of B_N and u_N .

Set

$$B_{N,n}(x) = \int_H B_N(e^{\frac{1}{n}A}x + y) N(0, Q_{\frac{1}{n}})(dy), \quad x \in H. \quad (33)$$

Then $B_{N,n}$ is of C^∞ class and all its derivatives are bounded. Moreover, $\|B_{N,n}\|_0 \leq \|B_N\|_0$, $n \geq 1$. It is easy to see that, possibly passing to a subsequence,

$$B_{N,n} \rightarrow B_N, \quad \mu - a.e.. \quad (34)$$

(indeed $B_{N,n} \rightarrow B_N$ in $L^2(H, \mu; H)$; this result can be first checked for continuous and bounded B). Now we denote by $u_{N,n}^{(i)}$ the solution of the equation

$$\lambda u_{N,n}^{(i)} - Lu_{N,n}^{(i)} - \langle B_{N,n}, Du_{N,n}^{(i)} \rangle = B_{N,n}^{(i)}, \quad (35)$$

where $B_{N,n}^{(i)} = \langle B_{N,n}, e_i \rangle$. By Lemma 8 we have, possibly passing to a subsequence, for any $x \in H$,

$$\begin{aligned} \lim_{n \rightarrow \infty} u_{N,n}^{(i)}(x) &= u_N^{(i)}(x), \quad \lim_{n \rightarrow \infty} Du_{N,n}^{(i)}(x) = Du_N^{(i)}(x), \\ \sup_{n \geq 1} \|u_{N,n}^{(i)}\|_{C_b^1(H)} &= C_{i,N} < \infty, \end{aligned} \quad (36)$$

where $u_N^{(i)}$ is the solution of (31).

Step 2 Approximation of X_t .

For any $m \geq i$ we set $X_m = (X_{m,t})$, $X_{m,t} := \pi_m X_t$, where $\pi_m = \sum_{j=1}^m e_j \otimes e_j$. Then we have

$$X_{m,t} = \pi_m x + \int_0^t A_m X_s ds + \int_0^t \pi_m B(X_s) ds + \pi_m W_t, \quad (37)$$

where $A_m = \pi_m A$.

Now we denote by $u_{N,n,m}^{(i)}$ the solution of the equation

$$\lambda u_{N,n,m}^{(i)} - L u_{N,n,m}^{(i)} - \langle \pi_m B_{N,n} \circ \pi_m, Du_{N,n,m}^{(i)} \rangle = B_{N,n}^{(i)} \circ \pi_m, \quad (38)$$

where $(B_{N,n} \circ \pi_m)(x) = B_{N,n}(\pi_m x)$, $x \in H$. Since only a finite number of variables is involved, we have, equivalently,

$$\lambda u_{N,n,m}^{(i)} - L^{(m)} u_{N,n,m}^{(i)} - \langle \pi_m B_{N,n} \circ \pi_m, Du_{N,n,m}^{(i)} \rangle = B_{N,n}^{(i)} \circ \pi_m,$$

with

$$L^{(m)} \varphi(x) = \frac{1}{2} \text{Tr} [\pi_m D^2 \varphi(x)] + \langle A_m x, D\varphi(x) \rangle. \quad (39)$$

Moreover, since $u_{N,n,m}^{(i)}$ depends only on the first m variables, we have

$$u_{N,n,m}^{(i)}(\pi_m y) = u_{N,n,m}^{(i)}(y), \quad y \in H. \quad (40)$$

Applying the finite-dimensional Itô formula to $u_{N,n,m}^{(i)}(X_{m,t}) = u_{N,n,m}^{(i)}(X_t)$ with the stopping time τ_N^X yields

$$\begin{aligned} & u_{N,n,m}^{(i)}(X_{m,t \wedge \tau_N^X}) - u_{N,n,m}^{(i)}(\pi_m x) = \\ & \int_0^{t \wedge \tau_N^X} \left(\frac{1}{2} \text{Tr} [D^2 u_{N,n,m}^{(i)}(X_{m,s})] + \langle Du_{N,n,m}^{(i)}(X_{m,s}), A_m X_s + \pi_m B(X_s) \rangle \right) ds \\ & + \int_0^{t \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_{m,s}), \pi_m dW_s \rangle. \end{aligned} \quad (41)$$

On the other hand, by (38) we have

$$\begin{aligned} & \lambda u_{N,n,m}^{(i)}(X_{m,t}) - \frac{1}{2} \text{Tr} [D^2 u_{N,n,m}^{(i)}(X_{m,t})] \\ & - \langle Du_{N,n,m}^{(i)}(X_{m,t}), A_m X_{m,t} + \pi_m B_{N,n}(X_{m,t}) \rangle = B_{N,n}^{(i)}(X_{m,t}). \end{aligned}$$

Let us fix $r \in]0, t]$. This will be useful in Step 3 of the proof to apply the Girsanov theorem (see in particular (49)).

Comparing with (41) and using (40) we find

$$u_{N,n,m}^{(i)}(X_{t \wedge \tau_N^X}) - u_{N,n,m}^{(i)}(X_{r \wedge \tau_N^X}) \quad (42)$$

$$\begin{aligned}
&= \lambda \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} u_{N,n,m}^{(i)}(X_s) ds - \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} B_{N,n}^{(i)}(X_{m,s}) ds \\
&+ \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_s), \pi_m(B(X_s) - B_{N,n}(X_{m,s})) \rangle ds + \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_s), dW_s \rangle.
\end{aligned}$$

Possibly passing to a subsequence, and taking the limit in probability as $m \rightarrow \infty$ (with respect to \mathbb{P}), we arrive at

$$\begin{aligned}
&u_{N,n}^{(i)}(X_{t \wedge \tau_N^X}) - u_{N,n}^{(i)}(X_{r \wedge \tau_N^X}) \tag{43} \\
&= \lambda \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} u_{N,n}^{(i)}(X_s) ds - \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} B_{N,n}^{(i)}(X_s) ds \\
&+ \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s), (B(X_s) - B_{N,n}(X_s)) \rangle ds + \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s), dW_s \rangle.
\end{aligned}$$

Let us justify this assertion.

First note that in equation (38) we have the drift term $\pi_m B_{N,n} \circ \pi_m$ which converges pointwise to $B_{N,n}$ and $B_{N,n}^{(i)} \circ \pi_m$ which converges pointwise to $B_{N,n}^{(i)}$ as $m \rightarrow \infty$. Since such functions are also uniformly bounded, we can apply Proposition 7 and Lemma 8 and obtain that, possibly passing to a subsequence (recall that n is fixed),

$$\begin{aligned}
\lim_{m \rightarrow \infty} u_{N,n,m}^{(i)}(x) &= u_{N,n}^{(i)}(x), \quad \lim_{m \rightarrow \infty} Du_{N,n,m}^{(i)}(x) = Du_{N,n}^{(i)}(x), \quad x \in H, \tag{44} \\
\sup_{m \geq 1} \|u_{N,n,m}^{(i)}\|_{C_b^1(H)} &= C_i^N < \infty.
\end{aligned}$$

We only consider convergence of the two most involved terms in (42).

We first treat convergence in $L^2(\Omega)$ of the stochastic integral. Recall that

$$\int_0^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s), dW_s \rangle = \int_0^t \mathbf{1}_{\{s \leq t \wedge \tau_N^X\}} \langle Du_{N,n}^{(i)}(X_s), dW_s \rangle;$$

by the isometry formula and (44) we get

$$\mathbb{E} \left| \int_0^{t \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_s) - Du_{N,n}^{(i)}(X_s), dW_s \rangle \right|^2 \rightarrow 0 \tag{45}$$

as $m \rightarrow \infty$. Note that we have used that $\lim_{m \rightarrow \infty} Du_{N,n,m}^{(i)}(x) = Du_{N,n}^{(i)}(x)$, for any $x \in H$ (not only for μ -a.e $x \in H$). In a similar way we get

$$\mathbb{E} \left| \int_0^{r \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_s) - Du_{N,n}^{(i)}(X_s), dW_s \rangle \right|^2 \rightarrow 0$$

as $m \rightarrow \infty$. This shows that $\int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n,m}^{(i)}(X_s), dW_s \rangle$ converges to $\int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s), dW_s \rangle$ in $L^2(\Omega)$ as $m \rightarrow \infty$. To show that, \mathbb{P} -a.s.,

$$\begin{aligned} \lim_{m \rightarrow \infty} \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} & \left| \langle Du_{N,n,m}^{(i)}(X_s), \pi_m(B(X_s) - B_{N,n}(X_{m,s})) \rangle \right. \\ & \left. - \langle Du_{N,n}^{(i)}(X_s), (B(X_s) - B_{N,n}(X_s)) \rangle \right| ds = 0, \end{aligned} \quad (46)$$

it is enough to prove that $\lim_{m \rightarrow \infty} H_m + K_m = 0$, where

$$H_m = \int_0^{t \wedge \tau_N^X} \left| \langle Du_{N,n,m}^{(i)}(X_s) - Du_{N,n}^{(i)}(X_s), \pi_m(B(X_s) - B_{N,n}(X_{m,s})) \rangle \right| ds$$

and

$$K_m = \int_0^{t \wedge \tau_N^X} \left| \langle Du_{N,n}^{(i)}(X_s), [\pi_m B(X_s) - B(X_s)] + [B_{N,n}(X_s) - \pi_m B_{N,n}(X_{m,s})] \rangle \right| ds.$$

By using (44) we easily get the assertion.

Step 3 A convergence result involving stopping times.

In order to pass to the limit in probability as $n \rightarrow \infty$ in (43) we recall formula (36) and argue as before. The only difficult term is

$$\int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s), (B(X_s) - B_{N,n}(X_s)) \rangle ds = J_n + I_n,$$

$$\text{where } J_n = \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s) - Du_N^{(i)}(X_s), (B(X_s) - B_{N,n}(X_s)) \rangle ds,$$

$$I_n = \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_N^{(i)}(X_s), (B_N(X_s) - B_{N,n}(X_s)) \rangle ds$$

(using that $s \leq t \wedge \tau_N^X$). As for J_n we have

$$J_n = \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_{N,n}^{(i)}(X_s) - Du_N^{(i)}(X_s), (B_N(X_s) - B_{N,n}(X_s)) \rangle ds,$$

and so $|J_n| \leq 2\|B_N\|_0 \int_0^t |Du_{N,n}^{(i)}(X_s) - Du_N^{(i)}(X_s)| ds \rightarrow 0$, \mathbb{P} -a.s., as $n \rightarrow \infty$, by Lemma 8.

Let us consider I_n . We define an auxiliary process $\hat{X}^N = (\hat{X}_t^N)$ as follows:

$$\hat{X}_t^N := e^{tA}x + \int_0^t e^{(t-s)A} B_N(X_{s \wedge \tau_N^X}) ds + \int_0^t e^{(t-s)A} dW_s, \quad t \geq 0, \quad (47)$$

Note that $X_{s \wedge \tau_N^X} = \hat{X}_{s \wedge \tau_N^X}^N$, $s \geq 0$, so that

$$\begin{aligned} |I_n| &= \left| \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_N^{(i)}(\hat{X}_s^N), (B_N(\hat{X}_s^N) - B_{N,n}(\hat{X}_s^N)) \rangle ds \right| \\ &\leq \|D^{(i)}u_N\|_0 \int_r^t |B_N(\hat{X}_s^N) - B_{N,n}(\hat{X}_s^N)| ds. \end{aligned} \quad (48)$$

Now we use the Girsanov theorem (see e.g. Appendix in [5]). Let $T > 0$. Since

$$\hat{X}_t^N := e^{tA}x + \int_0^t e^{(t-s)A} \hat{B}_s^N ds + \int_0^t e^{(t-s)A} dW_s, \quad t \geq 0,$$

where $\hat{B}_s^N = B_N(\hat{X}_{s \wedge \tau_N^X}^N)$, $s \geq 0$, is an adapted and bounded process, we have that

$$\tilde{W}_t^N := W_t + \int_0^t \hat{B}_s^N ds$$

is a cylindrical Wiener process on $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t \in [0, T]}, \tilde{\mathbb{P}}_N)$ where $\frac{d\tilde{\mathbb{P}}_N}{d\mathbb{P}} \Big|_{\mathcal{F}_T} = \rho_N$,

$$\rho_N = \exp \left(- \int_0^T \hat{B}_s^N dW_s - \frac{1}{2} \int_0^T |\hat{B}_s^N|^2 ds \right).$$

Hence $\hat{X}_t^N = e^{tA}x + \int_0^t e^{(t-s)A} d\tilde{W}_s^N$ is an OU process on $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t \in [0, T]}, \tilde{\mathbb{P}}_N)$.

Moreover, we know that the law of $(\hat{X}_t^N)_{t \in [0, T]}$ on $C([0, T]; H)$ is equivalent to the law of the OU process $Z(t, x)$ given in (16). In particular, all their transition probabilities are equivalent. Now under our assumptions the law of $Z(t, x)$ is equivalent to μ for all $t > 0$ and $x \in H$ (see Theorem 11.3 in [6]).

Let us come back to (48). Using that the law $\pi_t^N(x, \cdot)$ of \hat{X}_t^N is absolutely continuous with respect to μ , we obtain

$$\mathbb{E} \int_r^t |B_N(\hat{X}_s^N) - B_{N,n}(\hat{X}_s^N)| ds = \int_r^t ds \int_H |B_N(y) - B_{N,n}(y)| \frac{d\pi_s^N(x, \cdot)}{d\mu}(y) \mu(dy), \quad (49)$$

which tends to 0, as $n \rightarrow \infty$, by the dominated convergence theorem. Hence we have found that $I_n \rightarrow 0$ in $L^1(\Omega, \mathbb{P})$.

Up to now we have

$$\begin{aligned} &u_N^{(i)}(X_{t \wedge \tau_N^X}) - u_N^{(i)}(X_{r \wedge \tau_N^X}) \\ &= \lambda \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} u_N^{(i)}(X_s) ds - \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} B^{(i)}(X_s) ds + \int_{r \wedge \tau_N^X}^{t \wedge \tau_N^X} \langle Du_N^{(i)}(X_s), dW_s \rangle. \end{aligned}$$

Passing to the limit as $r \rightarrow 0^+$, since the trajectories of X are continuous, we finally get

$$\begin{aligned} & u_N^{(i)}(X_{t \wedge \tau_N^X}) - u_N^{(i)}(x) \\ &= \lambda \int_0^{t \wedge \tau_N^X} u_N^{(i)}(X_s) ds - \int_0^{t \wedge \tau_N^X} B^{(i)}(X_s) ds + \int_0^{t \wedge \tau_N^X} \langle Du_N^{(i)}(X_s), dW_s \rangle. \end{aligned} \quad (50)$$

Step 4 The final formula.

By (1) we deduce

$$dX_t^{(i)} = -\lambda_i X_t^{(i)} dt + B^{(i)}(X_t) dt + dW_t^{(i)}.$$

Inserting the expression for $B^{(i)}(X_t)$, which we get from this identity, into (50), we obtain

$$\begin{aligned} & u_N^{(i)}(X_{t \wedge \tau_N^X}) - u_N^{(i)}(x) \\ &= -X_{t \wedge \tau_N^X}^i + x^i + \lambda \int_0^{t \wedge \tau_N^X} u_N^{(i)}(X_s) ds - \lambda_i \int_0^{t \wedge \tau_N^X} X_s^{(i)} ds + \int_0^{t \wedge \tau_N^X} \langle Du_N^{(i)}(X_s), dW_s \rangle \\ & \quad + W_{t \wedge \tau_N^X}^i. \end{aligned}$$

By the variation of constants formula this is equivalent to

$$\begin{aligned} X_{t \wedge \tau_N^X}^{(i)} &= e^{-\lambda_i t \wedge \tau_N^X} \langle x, e_i \rangle + \lambda \int_0^{t \wedge \tau_N^X} e^{-\lambda_i(t \wedge \tau_N^X - s)} u_N^{(i)}(X_s) ds \\ & - \int_0^{t \wedge \tau_N^X} e^{-\lambda_i(t \wedge \tau_N^X - s)} du_N^{(i)}(X_s) + e^{-\lambda_i(t \wedge \tau_N^X)} \int_0^{t \wedge \tau_N^X} e^{\lambda_i s} [dW_s^{(i)} + \langle Du_N^{(i)}(X_s), dW_s \rangle]. \end{aligned}$$

This identity yields (32) on $\{t \leq \tau_N^X\}$. \square

The next lemma is similar to Lemma 9 in [5] and shows that $u_N(x) = \sum_{k \geq 1} u_N^{(k)}(x) e_k$ (where $u_N^{(k)}$ is as in (31)) is a well defined function which belongs to $C_b^1(H, H)$.

Lemma 11. *For λ sufficiently large, i.e., $\lambda \geq \tilde{\lambda}$, with $\tilde{\lambda} = \tilde{\lambda}(A, \|B_N\|_0) > 0$ there exists a unique $u_N = u_{\lambda, N} \in C_b^1(H, H)$ which solves*

$$u_N(x) = \int_0^\infty e^{-\lambda t} R_t (Du_N(\cdot) B_N(\cdot) + B_N(\cdot))(x) dt, \quad x \in H,$$

where R_t is the OU semigroup defined as in (14) and acting on H -valued functions. Moreover, we have the following assertions.

- (i) For any $h \in H$, $Du_N(\cdot)[h] \in C_b(H, H)$ and $\|Du_N(\cdot)[h]\|_0 \leq C_{0, \lambda, N} |h|$;
- (ii) for any $k \geq 1$, $\langle u_N(\cdot), e_k \rangle = u_N^{(k)}$, where $u_N^{(k)}$ is the solution defined in (31);
- (iii) There exists $c_3 = c_3(A, \|B_N\|_0) > 0$ such that, for any $\lambda \geq \tilde{\lambda}$, $u = u_\lambda$ satisfies

$$\|Du_N\|_0 \leq \frac{c_3}{\sqrt{\lambda}}. \quad (51)$$

3 Proof of Theorem 1

Let $X = (X_t)$ and $Y = (Y_t)$ be two weak mild solutions (see (8) and (9)) defined on the same filtered probability space (solutions with respect to the same cylindrical Wiener process W) starting at $x \in H$.

In the first part of the proof we will adapt the proof of Theorem 1 in [5] by introducing additional stopping times. The main difference with respect to [5] will appear in Proposition 13 which is needed to finish the proof.

For the time being, x is not specified. In Proposition 13 a restriction on x will emerge.

Note that by our hypothesis

$$B_N = B1_{B(0,N)} = B'1_{B(0,N)} = B'_N. \quad (52)$$

It follows that Kolmogorov equation (31) written with respect to the truncated drift B_N (with $B_N^{(i)}$ in the right-hand side) or with respect to B'_N (with $B_N^{(i)}$ in the right-hand side) is the same and gives the same solution $u_{N,\lambda}^{(i)}$.

It follows that both X and Y satisfy (32) on the event $\{t \leq \tau_N^X \wedge \tau_N^Y\}$.

Now we consider

$$u_N = u_{N,\lambda} : H \rightarrow H$$

be such that $u_N(x) = \sum_{i \geq 1} u_N^{(i)}(x)e_i$, $x \in H$, where $u_N^{(i)} = u_{N,\lambda}^{(i)}$ solve (31) for some λ large enough possibly depending on N .

Let us fix $T > 0$. By (32), taking into account (52), we have, for $t \in [0, T \wedge \tau_N^X \wedge \tau_N^Y]$, \mathbb{P} -a.s.,

$$\begin{aligned} X_t - Y_t &= u_N(Y_t) - u_N(X_t) + (\lambda - A) \int_0^t e^{(t-s)A} (u_N(X_s) - u_N(Y_s)) ds \\ &\quad + \int_0^t e^{(t-s)A} (Du_N(X_s) - Du_N(Y_s)) dW_s. \end{aligned}$$

Here and in the sequel we will drop the λ -dependence of u_N to simplify notation. However, at the end we will fix a value of λ large enough. By (51) we may assume that $\|Du_N\|_0 \leq 1/2$.

It follows that for $t \in [0, T \wedge \tau_N^X \wedge \tau_N^Y]$,

$$\begin{aligned} |X_t - Y_t| &\leq \frac{1}{2} |X_t - Y_t| + \left| (\lambda - A) \int_0^t e^{(t-s)A} (u_N(X_s) - u_N(Y_s)) ds \right| \\ &\quad + \left| \int_0^t e^{(t-s)A} (Du_N(X_s) - Du_N(Y_s)) dW_s \right|. \end{aligned}$$

Let η be a stopping time to be specified later and set

$$\tau = \eta \wedge T \wedge \tau_N^X \wedge \tau_N^Y. \quad (53)$$

Using that $1_{[0,\tau]}(t) = 1_{[0,\tau]}(t) \cdot 1_{[0,\tau]}(s)$, $0 \leq s \leq t \leq T$, we have (cf. page 187 in [6])

$$\begin{aligned} 1_{[0,\tau]}(t) |X_t - Y_t| &\leq C 1_{[0,\tau]}(t) \left| (\lambda - A) \int_0^t e^{(t-s)A} (u_N(X_s) - u_N(Y_s)) ds \right| \\ &+ C \left| 1_{[0,\tau]}(t) \int_0^t e^{(t-s)A} (Du_N(X_s) - Du_N(Y_s)) 1_{[0,\tau]}(s) dW_s \right|, \end{aligned}$$

where by C we denote any constant which may depend on the assumptions on A , B_N and T .

Writing $1_{[0,\tau]}(s) X_s = \tilde{X}_s$ and $1_{[0,\tau]}(s) Y_s = \tilde{Y}_s$, and, using the Burkholder-Davis-Gundy inequality with $q > 2$ which will be determined below, we obtain (recall that $\|\cdot\|$ is the Hilbert-Schmidt norm, cf. Chapter 4 in [6]) with $C = C_q$,

$$\begin{aligned} \mathbb{E} \left[\left| \tilde{X}_t - \tilde{Y}_t \right|^q \right] &\leq C \mathbb{E} \left[e^{\lambda qt} \left| (\lambda - A) \int_0^t e^{(t-s)A} e^{-\lambda s} (u_N(X_s) - u_N(Y_s)) 1_{[0,\tau]}(s) ds \right|^q \right] \\ &+ C \mathbb{E} \left[\left(\int_0^t 1_{[0,\tau]}(s) \left\| e^{(t-s)A} (Du_N(X_s) - Du_N(Y_s)) \right\|^2 ds \right)^{q/2} \right]. \end{aligned}$$

In the sequel we also introduce a parameter $\theta > 0$ and C_θ will denote suitable constants such that $C_\theta \rightarrow 0$ as $\theta \rightarrow +\infty$ (the constants may change from line to line). Similarly, we will indicate by $C(\lambda)$ suitable constants (possibly depending on N) such that $C(\lambda) \rightarrow 0$ as $\lambda \rightarrow +\infty$.

From the previous inequality we deduce, multiplying by $e^{-q\theta t}$, for any $\theta > 0$,

$$\begin{aligned} \mathbb{E} \left[e^{-q\theta t} \left| \tilde{X}_t - \tilde{Y}_t \right|^q \right] & \tag{54} \\ &\leq C \mathbb{E} \left[\left| (\lambda - A) \int_0^t e^{-\theta(t-s)} e^{(t-s)A} (u_N(X_s) - u_N(Y_s)) e^{-\theta s} 1_{[0,\tau]}(s) ds \right|^q \right] \\ &+ C \mathbb{E} \left[\left(\int_0^t e^{-2\theta(t-s)} \left\| e^{(t-s)A} (Du_N(X_s) - Du_N(Y_s)) \right\|^2 e^{-2\theta s} 1_{[0,\tau]}(s) ds \right)^{q/2} \right]. \end{aligned}$$

Now proceeding as in the proof of Theorem 7 of [5] we arrive at

$$\begin{aligned} &\int_0^T \mathbb{E} \left[e^{-q\theta t} \left| \tilde{X}_t - \tilde{Y}_t \right|^q \right] dt & \tag{55} \\ &\leq C(\lambda) \int_0^T e^{-q\theta s} \mathbb{E} |\tilde{X}_s - \tilde{Y}_s|^q ds + \tilde{C}_\theta \mathbb{E} \left[\Lambda_T \int_0^T e^{-q\theta s} |\tilde{X}_s - \tilde{Y}_s|^q ds \right] \end{aligned}$$

provided that $q \in (4, \infty)$, $\gamma = q/2$, $\theta \geq \lambda$ and

$$\begin{aligned} \Lambda_T &:= \int_0^T \int_0^t 1_{[0,\tau]}(s) \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^r)\|^2 \right)^\gamma dr ds dt, \\ &\text{where } Z_t^r = rX_t + (1-r)Y_t, \end{aligned} \tag{56}$$

and $D^2 u_N^{(n)}(x)$ is defined for μ -a.e. $x \in H$. The existence of $D^2 u_N^{(n)} \in L^p(\mu)$, $p \geq 2$, follows from Proposition 6 applied to equation (31) (see also Lemma 23 in [5]).

Since

$$\Lambda_T \leq T \cdot \int_0^{T \wedge \tau} \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^r)\|^2 \right)^\gamma dr ds$$

it is natural to define, for any $R > 0$, the stopping time

$$\bar{\tau}_R^{x,N} = \inf \left\{ t \geq 0 : \int_0^t \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^r)\|^2 \right)^\gamma dr ds \geq R \right\} \wedge T.$$

Take $\eta = \bar{\tau}_R^{x,N}$ in the previous expressions so that

$$\tau = \bar{\tau}_R^{x,N} \wedge \tau_N^X \wedge \tau_N^Y.$$

We get from (55),

$$\begin{aligned} & \int_0^T e^{-q\theta t} \mathbb{E} |\tilde{X}_t - \tilde{Y}_t|^q dt \\ & \leq C(\lambda) \int_0^T e^{-q\theta s} \mathbb{E} |\tilde{X}_s - \tilde{Y}_s|^q ds + \tilde{C}_\theta R \int_0^T e^{-q\theta s} \mathbb{E} |\tilde{X}_s - \tilde{Y}_s|^q ds. \end{aligned}$$

Now we fix λ large enough such that $C(\lambda) < 1/2$. For sufficiently large $\theta = \theta_R \geq \lambda$, depending on R and N , we have $\tilde{C}_\theta R < 1/2$ and so

$$\mathbb{E} \left[\int_0^T e^{-q\theta_R t} \mathbf{1}_{[0,\tau]}(t) |X_t - Y_t|^q dt \right] = \mathbb{E} \left[\int_0^\tau e^{-q\theta_R t} |X_t - Y_t|^q dt \right] = 0.$$

In other words, for every $R > 0$, $N \geq 1$, \mathbb{P} -a.s., $X = Y$ on $[0, \bar{\tau}_R^{x,N} \wedge \tau_N^X \wedge \tau_N^Y]$ (identically in t , not only a.e. in t , since X and Y are continuous processes).

If we prove that

$$\lim_{R \rightarrow \infty} \bar{\tau}_R^{x,N} = T \wedge \tau_N^X \wedge \tau_N^Y, \quad \mathbb{P} - a.s., \quad (57)$$

then we obtain that $X = Y$ on $[0, T \wedge \tau_N^X \wedge \tau_N^Y]$ and this finishes the proof.

The crucial assertion (57) follows by the next proposition.

Remark 12. Assertion (57) is a ‘‘local version’’ of Proposition 10 in [5]. Similarly to (47) also in the next proof we have to find an auxiliary process (see (58)) which allows to apply the Girsanov theorem.

Proposition 13. *Let $N \geq 1$ and $T > 0$ and suppose X and Y as in Theorem 1. For μ -a.e. $x \in H$, we have $\mathbb{P}(S_{T \wedge \tau_N^X \wedge \tau_N^Y}^x < \infty) = 1$, where*

$$S_t^x = S_t^{x,N} = \int_0^t \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^r)\|^2 \right)^\gamma dr ds, \quad t \in [0, T],$$

with $\gamma = q/2$ ($u_N(x) = \sum_{i \geq 1} u_N^{(i)}(x) e_i$, $x \in H$, where $u_N^{(i)} = u_{N,\lambda}^{(i)}$ solve (31)).

Proof. To prove the assertion we will show that, for μ -a.e. $x \in H$,

$$\mathbb{E}[S_{T \wedge \tau_N^X \wedge \tau_N^Y}^x] < +\infty.$$

In the first part of the proof, $x \in H$ is given, without restriction. Let us consider stopped processes

$$X_t^N = X_{t \wedge \tau_N^X \wedge \tau_N^Y}, \quad Y_t^N = Y_{t \wedge \tau_N^X \wedge \tau_N^Y},$$

and then we define an auxiliary process $(Z_t^{r,N})_{t \in [0, T]}$ as follows

$$Z_t^{r,N} := e^{tA}x + \int_0^t e^{(t-s)A} \bar{B}_s^{r,N} ds + \int_0^t e^{(t-s)A} dW_s \quad (58)$$

where (recall (10))

$$\bar{B}_s^{r,N} := [rB(X_s^N) + (1-r)B(Y_s^N)], \quad r \in [0, 1], \quad s \in [0, T].$$

Comparing Z^r (see (56)) and $Z^{r,N}$ we see that $Z_{s \wedge \tau_N^X \wedge \tau_N^Y}^r = Z_{s \wedge \tau_N^X \wedge \tau_N^Y}^{r,N}$, $s \in [0, T]$, $r \in [0, 1]$. Hence we have to prove

$$\mathbb{E}[S_{T \wedge \tau_N^X \wedge \tau_N^Y}^x] = \mathbb{E} \int_0^{T \wedge \tau_N^X \wedge \tau_N^Y} \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^{r,N})\|^2 \right)^\gamma dr ds < \infty.$$

This follows if we can show that

$$\mathbb{E} \int_0^T \int_0^1 \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(Z_s^{r,N})\|^2 \right)^\gamma dr ds < \infty. \quad (59)$$

We fix $N \geq 1$. To verify (59) we can follow the proof of Proposition 10 in [5]. We only indicate the small changes which are needed.

Define

$$\rho_{r,N} = \exp \left(- \int_0^t \bar{B}_s^{r,N} dW_s - \frac{1}{2} \int_0^t |\bar{B}_s^{r,N}|^2 ds \right).$$

We have, since $|\bar{B}_s^r| \leq \|B_N\|_0$, $r \in [0, 1]$, $s \geq 0$, \mathbb{P} -a.s.,

$$\mathbb{E} \left[\exp \left(k \int_0^T |\bar{B}_s^{r,N}|^2 ds \right) \right] \leq C_k < \infty, \quad (60)$$

for all $k \in \mathbb{R}$, independently of x and r , simply because B_N is bounded. Hence an infinite dimensional version of Girsanov's Theorem with respect to a cylindrical Wiener process (the proof of which is included in the Appendix of [5]) applies and gives us that

$$\tilde{W}_t^N := W_t + \int_0^t \bar{B}_s^{r,N} ds$$

is a cylindrical Wiener process on $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t \in [0, T]}, \tilde{\mathbb{P}}_{r, N})$ where $\frac{d\tilde{\mathbb{P}}_{r, N}}{d\mathbb{P}} \Big|_{\mathcal{F}_T} = \rho_{r, N}$. Hence

$$Z_t^{r, N} = e^{tA}x + \int_0^t e^{(t-s)A} d\tilde{W}_s^N$$

is an OU process on $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t \in [0, T]}, \tilde{\mathbb{P}}_{r, N})$. Continuing as in the proof of Proposition 10 in [5] with $u^{(n)}$ replaced by $u_N^{(n)}$, ρ_r replaced by $\rho_{r, N}$, Z_s^r replaced by $Z_s^{r, N}$, \bar{B}_s replaced by $\bar{B}_s^{r, N}$, we see that (59) holds if we prove that

$$\int_0^T \mathbb{E} \left[\left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(e^{sA}x + W_A(s))\|^2 \right)^{2\gamma} \right] ds < \infty, \quad (61)$$

where $W_A(t) = \int_0^t e^{(t-s)A} dW(s)$, $t \geq 0$. If μ_s^x denotes the law of $e^{sA}x + W_A(s)$, we have to prove that

$$\int_0^T \int_H \left(\sum_{n \geq 1} \frac{1}{\lambda_n^{1-\delta}} \|D^2 u_N^{(n)}(y)\|^2 \right)^{2\gamma} \mu_s^x(dy) ds < \infty. \quad (62)$$

This can be checked as in the mentioned proof (see in particular Steps 3 and 4 in that proof) *only for μ -a.e. $x \in H$* ; one has to replace B in the proof in [5] with our B_N . \square

Remark 14. As is easily checked in Theorem 1 the ball $B(0, N)$ can be replaced by any open bounded set in H .

Remark 15. According to Remark 11 in [5] our Theorem 1 provides an alternative approach to Veretennikov's uniqueness result in finite dimension. In this respect first note that Theorem 1 when $H = \mathbb{R}^d$ does not require to start from μ -a.e. initial conditions x , but works for any initial $x \in \mathbb{R}^d$.

Note that in finite dimension an SDE like $dX_t = b(X_t)dt + dW_t$ with b Borel and bounded is equivalent to $dX_t = -X_t dt + (b(X_t) + X_t)dt + dW_t$ which is in the form (1) with $A = -I$, and with a drift term $B(x) = b(x) + x$ which is completely covered by Theorem 1.

Recall that in this alternative approach to Veretennikov's result, basically the elliptic L^p -estimates with respect to Lebesgue measure used in [26] are replaced by elliptic $L^p(\mu)$ -estimates.

4 Existence of strong mild solutions

Here we will use our Theorem 1 to prove existence of strong mild solutions when B grows more than linearly. We will construct such solutions for μ -a.e. initial $x \in H$.

According to Chapter 1 in [21] (see also [18]) if $x \in H$ we say that equation (1) has a (global) *strong mild solution* if, for every filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ on

which there is defined a cylindrical (\mathcal{F}_t) -Wiener process W there exists an H -valued continuous (\mathcal{F}_t) -adapted process $X = (X_t) = (X_t)_{t \geq 0}$ such that

$$(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}, W, X)$$

is a weak mild solution.

Theorem 16. *Let us consider equation (1) and assume Hypothesis 1 and $B \in B_{b,loc}(H, H)$. Moreover, suppose that there exist $C > 0$, $p > 0$, such that*

$$\langle B(y+z), y \rangle \leq C(|y|^2 + e^{p|z|} + 1), \quad y, z \in H. \quad (63)$$

Then, for μ -a.e. $x \in H$ (where $\mu = N(0, -\frac{1}{2}A^{-1})$), equation (1) has a strong mild solution. Moreover, this solution is pathwise unique.

Remark 17. Condition (63) is a bit stronger than the classical one: $\langle B(y), y \rangle \leq C(|y|^2 + 1)$, $y \in H$, which is usually imposed in finite dimension to have non-explosion for SDEs with additive noise. We can not use such condition. Indeed for a given mild solution (X_t) we can not write the Itô formula for $|X_t|^2$ due to the fact that our noise is cylindrical.

To prove the result we will use our Theorem 1 together with a generalization of the Yamada-Watanabe theorem (see Theorem 2 in [21] and [18]) and some a-priori estimates on mild solutions (see Section 4.1).

Example 18. To introduce an example of a drift \tilde{B} which satisfies the assumptions of Theorem 16, we first consider a measurable function $g : \mathbb{R} \rightarrow \mathbb{R}_+$ such that $g(s) = 0$ if $s \geq 0$ and $0 \leq g(s) \leq Ce^{q|s|}$, $s < 0$, for some $q > 0$. It is easy to check that $g(s+r)r \leq C'(1 + |r|^2 + e^{p|s|})$, $s, r \in \mathbb{R}$. We define $\tilde{B} : H \rightarrow H$,

$$\tilde{B}(x) = \sum_{k \geq 1} \frac{g(x_k)}{k^2} e_k, \quad x \in H.$$

It is not difficult to verify that \tilde{B} satisfies the assumptions of the previous theorem. We can also add to our drift \tilde{B} one of the singular drifts considered in Section 4 of [5]; we will still obtain an admissible drift for our theorem.

4.1 An a-priori estimate

Here we prove an a-priori estimate for mild solutions to (1) under condition (63). For this purpose let us consider the OU process

$$Z_t = Z(t, x) = e^{tA}x + \int_0^t e^{(t-s)A} dW_s$$

which under our hypotheses has a continuous H -valued version. It satisfies

$$\langle Z_t, \varphi \rangle = \int_0^t \langle Z_s, A\varphi \rangle ds + \langle W_t, \varphi \rangle$$

for all $\varphi \in D(A)$. By Proposition 18 in [5] we deduce, in particular, that for any $p > 0$, $T > 0$

$$K_T := \mathbb{E} \left[\sup_{t \in [0, T]} e^{p|Z_t|} \right] < \infty.$$

Recall that under our hypotheses a weak mild solution to (1) can be defined, equivalently, as a tuple $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}, W, X)$, where $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ is a filtered probability space on which there is defined a cylindrical (\mathcal{F}_t) -Wiener process W and a continuous (\mathcal{F}_t) -adapted H -valued process $X = (X_t) = (X_t)_{t \geq 0}$ such that, \mathbb{P} -a.s.,

$$\langle X_t, \varphi \rangle = \langle x, \varphi \rangle + \int_0^t (\langle X_s, A\varphi \rangle + \langle B(X_s), \varphi \rangle) ds + \langle W_t, \varphi \rangle, \quad t \geq 0,$$

for all $\varphi \in D(A)$ (cf. Chapter 6 of [6]).

Theorem 19. *Assume Hypothesis 1, $B \in B_{b,loc}(H, H)$ and condition (63). Let $X = (X_t)_{t \geq 0}$ be a weak mild solution of equation (1) with $X_0 = x \in H$.*

There exists $C_p > 0$ (possibly depending on C and p given in (63)) such that

$$\mathbb{E} \left[\sup_{t \in [0, T]} |X_t|^2 \right] \leq e^{C_p T} (|x|^2 + K_T + 1), \quad T > 0. \quad (64)$$

Proof. The process $Y_t = X_t - Z_t$ satisfies

$$\langle Y_t, \varphi \rangle = \langle x, \varphi \rangle + \int_0^t (\langle Y_s, A\varphi \rangle + \langle B(Y_s + Z_s), \varphi \rangle) ds, \quad (65)$$

for all $\varphi \in D(A)$, and it has continuous trajectories in H . Let us consider (65) with $\varphi = e_k$ (see (13)) and set $Y_t^{(k)} = \langle Y_t, e_k \rangle$. Since $\frac{dY_t^{(k)}}{dt} \cdot Y_t^{(k)} \leq B^{(k)}(Y_t + Z_t)Y_t^{(k)}$, we find

$$\sum_{k \geq 1} \langle Y_t, e_k \rangle^2 \leq |x|^2 + 2 \int_0^t \langle B(Y_s + Z_s), Y_s \rangle ds.$$

Hence, by assumption (63), for $t \in [0, T]$,

$$|Y_t|^2 \leq |x|^2 + 2C \int_0^t (|Y_s|^2 + e^{p|Z_s|} + 1) ds,$$

and therefore, by the Gronwall lemma,

$$|Y_t|^2 \leq e^{CT} \left(|x|^2 + 2CT \left(\sup_{s \in [0, T]} e^{p|Z_s|} + 1 \right) \right).$$

This implies

$$\mathbb{E} \left[\sup_{t \in [0, T]} |Y_t|^2 \right] \leq e^{C_1 T} \left(|x|^2 + \mathbb{E} \left[\sup_{s \in [0, T]} e^{p|Z_s|} \right] + 1 \right).$$

Therefore,

$$\begin{aligned} \mathbb{E} \left[\sup_{t \in [0, T]} |X_t|^2 \right] &\leq 2e^{C_1 T} \left(|x|^2 + \mathbb{E} \left[\sup_{s \in [0, T]} e^{p|Z_s|} \right] + 1 \right) + 2\mathbb{E} \left[\sup_{t \in [0, T]} |Z_t|^2 \right] \\ &\leq e^{C_p T} \left(|x|^2 + \mathbb{E} \left[\sup_{s \in [0, T]} e^{p|Z_s|} \right] + 1 \right). \end{aligned}$$

□

4.2 Proof of Theorem 16

By Theorem 1 we only have to prove existence of strong solution for μ -a.e. $x \in H$.

We will again consider truncated bounded drifts $B_N = B \mathbf{1}_{B(0, N)}$, $N \geq 1$.

By the main result in [5] there exists a Borel set $\tilde{G} \subset H$ with $\mu(\tilde{G}) = 1$ such that for any $x \in \tilde{G}$ we have pathwise uniqueness for each stochastic equation

$$dX_t = (AX_t + B_N(X_t))dt + dW_t, \quad X_0 = x \in H, \quad (66)$$

$N \geq 1$. Let $x \in \tilde{G}$. By the Girsanov theorem (see Appendix in [5]) there exists (a unique in law) weak mild solution $X_N = (X_N(t))_{t \geq 0}$ for each stochastic equation (66).

Therefore we can apply a generalization of the Yamada-Watanabe theorem (see Theorem 2 in [21] and [18]) to (66) when $x \in \tilde{G}$.

Let us fix any filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ on which there is defined a cylindrical (\mathcal{F}_t) -Wiener process W . By the Yamada-Watanabe theorem, for any $N \geq 1$, on the fixed filtered probability space above there exists a (unique) strong mild solution X_N to (66). Moreover, since

$$B_N(x) = B_{N+k}(x), \quad x \in B(0, N),$$

$k \geq 1$, we have by Theorem 1 that, \mathbb{P} -a.s.,

$$\tau_N := \tau_N^{X_N} = \tau_N^{X_{N+k}}, \quad k \geq 1, \quad N \geq 1,$$

and $X_N(t \wedge \tau_N) = X_{N+k}(t \wedge \tau_N)$, $t \geq 0$.

It is enough to construct the strong solution X to (1) on $[0, T]$ for a fixed $T > 0$. We define an H -valued stochastic process X on $\Omega' = \cup_{N \geq 1} \{\tau_N > T\}$ as

$$X(t)(\omega) := X_N(t)(\omega), \quad t \in [0, T],$$

if $\omega \in \{\tau_N > T\}$ (we set $X_t(\omega) = 0$ if $\omega \notin \Omega'$, $t \in [0, T]$). Then $X(t)$ is well defined.

It is not difficult to prove that X is a strong mild solution on $[0, T]$ if we show that

$$\lim_{N \rightarrow \infty} \mathbb{P}(\tau_N > T) = 1 \quad (67)$$

(this will imply that $\mathbb{P}(\Omega') = 1$). To verify (67) we will apply Theorem 19. Note that each B_N satisfies

$$\langle B_N(y+z), y \rangle \leq C(|y|^2 + e^{p|z|} + 1), \quad y, z \in H,$$

with the same constants C and p of (63). By Theorem 19 we obtain

$$\mathbb{E} \left[\sup_{t \in [0, T]} |X_N(t)|^2 \right] \leq e^{C_p T} (|x|^2 + K_T + 1), \quad (68)$$

with C_p independent of $N \geq 1$. Since

$$\mathbb{P}(\tau_N \leq T) = \mathbb{P} \left(\sup_{t \in [0, T]} |X_N(t)| \geq N \right)$$

by (68) and the Chebychev inequality we easily get assertion (67) from (68) and this completes the proof.

4.3 Existence and uniqueness of local mild solutions

Finally, let us discuss a possible extension of our result to the case when the drift term B only belongs to $B_{b,loc}(H, H)$, without requiring hypothesis (63).

We need the concept of local solution (see, for instance, [1] for some additional facts about local solutions). Let $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ be a filtered probability space. A stopping time $\tau : \Omega \rightarrow [0, +\infty]$ is called *accessible* if there exists a sequence of stopping times $(\tau_n)_{n \in \mathbb{N}}$ such that $\mathbb{P}(\tau_n < \tau) = 1$ and $\mathbb{P}(\lim_{n \rightarrow \infty} \tau_n = \tau) = 1$. The previous sequence (τ_n) is called an approximating sequence of τ .

Notice that, if τ_1 and τ_2 are accessible stopping times, then also $\tau = \tau_1 \wedge \tau_2$ is an accessible stopping time.

Let τ be an accessible stopping time and consider $[0, \tau) \times \Omega = \{(t, \omega) \in [0, +\infty) \times \Omega : 0 \leq t < \tau(\omega)\}$. An H -valued stochastic process X defined on $[0, \tau)$ (i.e., $X : [0, \tau) \times \Omega \rightarrow H$) is called (\mathcal{F}_t) -adapted if $X_t(\cdot) : \{t < \tau\} \rightarrow H$ is \mathcal{F}_t -measurable, for any $t \geq 0$ (on $\{t < \tau\}$ we consider the restricted σ -algebra $\{A \cap \{t < \tau\}\}_{A \in \mathcal{F}_t}$); moreover, it is called continuous if trajectories are continuous on $[0, \tau)$, \mathbb{P} -a.s.. Note that X is (\mathcal{F}_t) -adapted if and only if the process $\tilde{X} = (\tilde{X}_t)_{t \geq 0}$,

$$\tilde{X}_t = X_t 1_{\{\tau > t\}} + 0 \cdot 1_{\{\tau \leq t\}} \quad (69)$$

is (\mathcal{F}_t) -adapted.

Definition 2. Let $x \in H$. We call local weak mild solution to (1) a tuple $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}, W, X, \tau)$, where $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ is a filtered probability space on which it is defined a cylindrical (\mathcal{F}_t) -Wiener process W , an accessible stopping time τ and a continuous (\mathcal{F}_t) -adapted H -valued process X defined on $[0, \tau)$ such that, there exists an approximating sequence (τ_n) of τ for which, \mathbb{P} -a.s., on $\{t \leq \tau_n\}$

$$X_t = e^{tA}x + \int_0^t e^{(t-s)A}B(X_s) ds + \int_0^t e^{(t-s)A}dW_s,$$

for all $n \in \mathbb{N}$ and $t \geq 0$.

A local weak mild solution X which is $(\bar{\mathcal{F}}_t^W)$ -adapted (here $(\bar{\mathcal{F}}_t^W)$ denotes the completed natural filtration of the cylindrical process W) and such that τ is an $(\bar{\mathcal{F}}_t^W)$ -stopping time is called a local strong mild solution.

Theorem 20. Assuming Hypothesis 1 and $B \in B_{b,loc}(H, H)$, existence of local strong mild solutions holds for μ -a.e. initial condition $x \in H$. Moreover, for μ -a.e. $x \in H$, if X and Y are two local weak mild solutions on the same $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ with the same cylindrical Wiener process W , defined on $[0, \tau^X)$ and $[0, \tau^Y)$ respectively (τ^X and τ^Y are accessible stopping times as in Definition 2), then, \mathbb{P} -a.s., $X = Y$ on $[0, \tau)$, where $\tau = \tau^X \wedge \tau^Y$.

Proof. Let us sketch some of the details of the proof.

From the first part of the proof of Theorem 16 (see also Theorem 1 in [5]), given a priori a filtered probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ and a cylindrical (\mathcal{F}_t) -Wiener process W , we have the existence of a local strong mild solution X_N on $[0, \tau_N)$, for every $N \in \mathbb{N}$. Note that each τ_N is accessible since as approximating sequence we may take $\tau_{N,n} = \inf\{t \geq 0 : X_N(t) \notin B(0, N - \frac{1}{n})\} \wedge n$, for $n \geq 1$.

Thus, taking $\tau = \sup_{N \in \mathbb{N}} \tau_N$, τ is accessible and we have a local strong mild solution on $[0, \tau)$.

In order to prove uniqueness, we first note that if X is a local weak mild solution defined on $[0, \tau^X)$ and (τ_n^X) is an approximating sequence of τ^X as in the definition of solution, then assertion (32) of Theorem 10 holds, for any $t > 0$, $n \geq 1$, \mathbb{P} -a.s., on the event $\{t \leq \tau_{n,N}^X\}$, where $\tau_{n,N}^X$ is the stopping time

$$\tau_{n,N}^X = \inf\{t \in [0, \tau^X) : X_t \notin B(0, N)\} \wedge \tau_n^X$$

($\tau_{n,N}^X = \tau_n^X$ if the set is empty; to show that $\tau_{n,N}^X$ is a stopping time, note that $\tau_{n,N}^X = \inf\{t \geq 0 : \tilde{X}_t \notin B(0, N)\} \wedge \tau_n^X$, where \tilde{X}_t is defined in (69) with τ replaced by τ^X).

Indeed, one can repeat the arguments in the proof of Theorem 10 with the same functions u_N , replacing τ_N^X with $\tau_{n,N}^X$.

Now let X and Y be two local weak mild solutions as in the second part of the theorem.

If (τ_n^X) and (τ_n^Y) are, respectively, approximating sequences of τ^X and τ^Y as in the definition of solution, then in order to prove uniqueness, it is enough to consider $\sigma_n = \tau_n^X \wedge \tau_n^Y$ and check that, \mathbb{P} -a.s.,

$$X = Y \text{ on } [0, \sigma_n], \quad n \geq 1. \quad (70)$$

Let us fix $n \geq 1$. We can adapt the proof of Theorem 1, arguing on the interval $[0, \eta \wedge T \wedge \tau_{n,N}^X \wedge \tau_{n,N}^Y]$ (cf. (53)). We finally get that $X = Y$ on $[0, T \wedge \sigma_n]$, $T > 0$, and this gives the assertion. \square

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