# AN ALTERNATIVE PROOF OF WELL-POSEDNESS OF STOCHASTIC EVOLUTION EQUATIONS IN THE VARIATIONAL SETTING 

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#### Abstract

We present a new proof of well-posedness of stochastic evolution equations in variational form, relying solely on a (nonlinear) infinite-dimensional approximation procedure rather than on classical finite-dimensional projection arguments of Galerkin type.


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Key words: stochastic evolution equations, variational approach, monotone operators.

## 1. INTRODUCTION

Let us consider a stochastic evolution equation of the type

$$
\begin{equation*}
d X(t)+A(t, X(t)) d t=B(t, X(t)) d W(t), \quad X(0)=X_{0} \tag{1.1}
\end{equation*}
$$

in the so-called variational setting, i.e. where $A$ is a random time-dependent nonlinear maximal monotone operator from a reflexive Banach space $V$ to its dual $V^{\prime}$, with $V$ densely and continuously embedded in a Hilbert space $H$. Moreover, $W$ is a cylindrical Wiener process (possibly defined on a further separable Hilbert space), and $B$ is a random time-dependent map with values in a suitable space of Hilbert-Schmidt operators. Precise assumptions on the data of the problem are given in $\S 2$ below.

This class of equations was introduced and studied by Pardoux in [11], extending to the stochastic setting the classical well-posedness results by Lions (see, e.g., [6]) for equations without noise. More precisely, in [11] the operator $A$ is time-dependent but non-random and the pair $(A, B)$ needs to satisfy coercivity and boundedness assumptions, complemented by a local Lipschitz continuity condition on $B$. The general case where $A$ can be random was considered by Krylov and Rozovskiĭ in [5], who also showed that the local Lipschitz
continuity on $B$ is not needed (see also [7, 12] for comprehensive treatments and recent developments). Under coercivity and boundedness assumptions on $A$ and a Lipschitz continuity assumption on $B$ (a set of hypotheses to which we shall refer to as disjoint assumptions), Pardoux [11, Chapter 3, §1] proved well-posedness of (1.1) by a clever, but at the same time natural extension of the deterministic theory, employing an infinite-dimensional argument based on Picard iterations in suitable spaces of processes. As a second step, assuming that the pair $(A, B)$ satisfies a joint coercivity and boundedness assumption and that $B$ is locally Lipschitz continuous, well-posedness for (1.1) is proved by a different method, i.e. by finite-dimensional approximations of Galerkin type (see [11, Chapter 3, §3]). As mentioned above, the joint assumption is shown to imply well-posedness without any local Lipschitz continuity condition in $B$ in [5], again using finite-dimensional approximations.

Our goal is to show that existence of solutions under the general joint assumptions on $(A, B)$, as in [5], can be obtained relying only on infinitedimensional arguments, i.e. in the same spirit of the approach adopted in the first part of [11]. Moreover, since the existence proof in [5] relies on quite advanced results for finite-dimensional stochastic differential equations, the alternative proof provided here could also be seen as a simpler proof. The main idea is, roughly speaking, to regularize the operator $B$ through the resolvent of the operator $A$. The corresponding regularized problem is then shown to satisfy the stronger disjoint hypotheses, so that it can be solved, as in [11, Chapter 3, § 1], using infinite-dimensional techniques only. Uniform estimates on the solutions to the regularized equations are then established, which allow to pass to the limit obtaining a solution to the original problem.

Even though the use of Yosida approximations is a standard tool in the field of nonlinear deterministic and stochastic equations (see, e.g., $[1,2,3,8]$ for just a few examples among an enormous literature), it seems that the type of approximation introduced here is not found elsewhere. Let us also mention that another approach to stochastic equations in variational form, namely by reduction to the deterministic case, is developed in [1, §4.4], where, however, only the case of additive noise is considered. Moreover, well-posedness for certain classes of stochastic equations with multi-valued nonlinear drift term is obtained in [8] using the results in $[5,11]$ as starting point. The main point in [8] is to regularize the multi-valued drift coefficient by its Yosida approximation, thus obtaining a family of well-posed equations with singlevalued drift satisfying the assumptions of the classical variational framework, and to show that the corresponding approximate solutions converge, under appropriate assumptions, among which the Lipschitz continuity of the diffusion coefficient, to a process solving the original equation.

## 2. SETTING AND MAIN RESULTS

Throughout the paper, $\left(\Omega, \mathscr{F},(\mathscr{F} t)_{t \in[0, T]}, \mathbb{P}\right)$ stands for a filtered probability space satisfying the so-called usual conditions, where $T>0$ is a fixed final time, on which all random elements will be defined. Equality of processes is always meant in the sense of indistinguishability, unless otherwise stated. Moreover, $U$ is a separable Hilbert space and $W$ is a cylindrical Wiener process on it; $H$ is a separable Hilbert space identified with its dual, and $V$ is a separable reflexive Banach space continuously and densely embedded in $H$, so that, denoting the (topological) dual of $V$ by $V^{\prime}, V \hookrightarrow H \hookrightarrow V^{\prime}$ is a Gelfand triple. The scalar product and norm of $H$ will be denoted by $\langle\cdot, \cdot\rangle$ and $\|\cdot\|$, respectively, while the norms of all other Banach spaces will be indicated by subscripts. Since the duality form between $V$ and $V^{\prime}$ agrees with the scalar product of $H$ in the usual sense, we shall denote the former by $\langle\cdot, \cdot\rangle$ as well. If $E_{1}$ and $E_{2}$ are Hilbert spaces, the space of Hilbert-Schmidt operators from $E_{1}$ to $E_{2}$ will be denoted by $\mathscr{L}^{2}\left(E_{1} ; E_{2}\right)$.

The following assumptions will be in force throughout the paper.
(I) The operator $A: \Omega \times[0, T] \times V \rightarrow V^{\prime}$ is progressively measurable and hemicontinuous, i.e., for every $x \in V$ the $V^{\prime}$-valued process $A(\cdot, \cdot, x)$ is progressively measurable and the map

$$
\mathbb{R} \ni r \longmapsto\langle A(\omega, t, x+r y), z\rangle
$$

is continuous for every $\omega \in \Omega, t \in[0, T]$, and $x, y, z \in V$.
(II) The operator $B: \Omega \times[0, T] \times V \rightarrow \mathscr{L}^{2}(U ; H)$ is progressively measurable, i.e., for every $x \in V$ the $\mathscr{L}^{2}(U ; H)$-valued process $B(\cdot, \cdot, x)$ is progressively measurable.
(III) There exist constants $\left.c_{1}>0, c_{2} \geqslant 0, p \in\right] 1,+\infty$ [ and an adapted process $f \in L^{1}(\Omega \times(0, T))$ such that

$$
\begin{gathered}
\langle A(\omega, t, x)-A(\omega, t, y), x-y\rangle-\frac{1}{2}\|B(\omega, t, x)-B(\omega, t, y)\|_{\mathscr{L}^{2}(U ; H)}^{2} \\
\geqslant-c_{2}\|x-y\|^{2} \\
\langle A(\omega, t, x), x\rangle-\frac{1}{2}\|B(\omega, t, x)\|_{\mathscr{L}^{2}(U ; H)}^{2} \geqslant c_{1}\|x\|_{V}^{p}-c_{2}\|x\|^{2}-f(\omega, t)
\end{gathered}
$$

for every $x, y \in V, t \in[0, T]$ and $\omega \in \Omega$.
(IV) There exist a constant $C>0$ and an adapted process $g \in L^{1}(\Omega \times(0, T))$ such that, setting $q:=p /(p-1)$,

$$
\|A(\omega, t, x)\|_{V^{\prime}}^{q} \leqslant C\|x\|_{V}^{p}+g(\omega, t)
$$

for every $x \in V, t \in[0, T]$ and $\omega \in \Omega$.
(V) $X_{0} \in L^{2}\left(\Omega, \mathscr{F}_{0} ; H\right)$.

We can give now the definition of strong solution for the equation.
Definition 2.1. A strong solution to (1.1) is a $V$-valued progressively measurable process $X$ such that

$$
\begin{gathered}
X \in L^{0}(\Omega ; C([0, T] ; H)) \cap L^{0}\left(\Omega ; L^{p}(0, T ; V)\right) \\
B(\cdot, \cdot, X) \in L^{0}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right),
\end{gathered}
$$

and

$$
X+\int_{0} A(s, X(s)) d s=X_{0}+\int_{0} B(s, X(s)) d W(s)
$$

as an identity in the sense of indistinguishable $V^{\prime}$-valued processes.
The classical well-posedness result $[5,11]$ for equation (1.1) is as follows.
Theorem 2.2. There exists a unique strong solution $X$ to (1.1). Moreover,
and the solution map

$$
X \in L^{2}(\Omega ; C([0, T] ; H)) \cap L^{p}\left(\Omega ; L^{p}(0, T ; V)\right)
$$

$$
\begin{aligned}
L^{2}\left(\Omega, \mathscr{F}_{0} ; H\right) & \longrightarrow C\left([0, T] ; L^{2}(\Omega ; H)\right) \\
X_{0} & \longmapsto X
\end{aligned}
$$

is Lipschitz continuous.
As discussed above, in the next section we show that Theorem 2.2 can be proved relying only on infinite-dimensional arguments.

Remark 2.3. The proof of uniqueness of strong solution crucially relies on an Itô formula for the square of the $H$-norm. In [11] such formula is obtained under the assumption that an operator $C: V \rightarrow V^{\prime}$ exists satisfying monotonicity, coercivity and boundedness conditions (see [11, p. 57]). These conditions coincide with those assumed here on $A(\omega, t)$ for all $(\omega, t) \in \Omega \times[0, T]$, hence are automatically verified. In general, the formula remains valid, without any connection to a specific stochastic equation, if the duality map $J: V \rightarrow V^{\prime}$ is single-valued, which is the case if, e.g., $V^{\prime}$ is a strictly convex Banach space. Such an assumption is always satisfied in all applications to SPDEs we know of. On the other hand, the Itô formula in [5] does not require any "geometric" assumption on the Banach space $V$, but its proof is rather involved and relies on finite-dimensional projections. A simple proof of Itô's formula for the square of the $H$-norm in the variational setting, which relies just on infinite-dimensional arguments, is available, to the best of our knowledge, only in the case where $V$ is a Hilbert space (see [4], as well as [9]).

## 3. PROOF OF THEOREM 2.2

With $\omega \in \Omega$ and $t \in[0, T]$ arbitrary but fixed, assumptions (I) and (III) imply that $\tilde{A}:=A(\omega, t, \cdot)+c_{2} I: V \rightarrow V^{\prime}$ is maximal monotone. Let us show that the part of $\tilde{A}$ in $H$, denoted by $\tilde{A}_{H}$, is maximal monotone (as an operator in $H$ ) with domain $\mathrm{D}(\tilde{A}):=\{x \in V: A(\omega, t, x) \in H\}$. It suffices to show that, for any $y \in H$, the equation

$$
x+\tilde{A} x=y
$$

admits a solution $x \in \mathrm{D}(\tilde{A})$. Since $\tilde{A}$ is coercive on $V$ by assumption (III), it follows by maximal monotonicity that the equation admits a (unique) solution $x \in V$. This obviously implies $\tilde{A} x \in H$, i.e. $x \in \mathrm{D}(\tilde{A})$. We have hence shown that $\tilde{A}_{H}$ is a maximal monotone operator on $H$.

For every $\lambda>0$ we define the resolvent operator

$$
J_{\lambda}: \Omega \times[0, T] \times H \longrightarrow V
$$

of $\tilde{A}$ setting, for every $x \in H, t \in[0, T]$ and $\omega \in \Omega$,

$$
J_{\lambda}(\omega, t, x)+\lambda \tilde{A}\left(\omega, t, J_{\lambda}(\omega, t, x)\right)=x .
$$

The maximal monotonicity of $\tilde{A}_{H}$ implies that, for every $(\omega, t) \in \Omega \times[0, T]$, $J_{\lambda}(\omega, t, \cdot)$ is a contraction and converges pointwise to the identity map of $H$ as $\lambda \rightarrow 0$.

Moreover, we define the Yosida approximation of $\tilde{A}$ as the map

$$
\begin{aligned}
\tilde{A}_{\lambda}: \Omega \times[0, T] \times H & \longrightarrow H \\
(\omega, t, x) & \longmapsto \tilde{A}\left(\omega, t, J_{\lambda}(\omega, t, x)\right) .
\end{aligned}
$$

It follows by the contraction property of $J_{\lambda}$ that, for every $(\omega, t) \in \Omega \times[0, T]$, $\tilde{A}_{\lambda}(\omega, t, \cdot)$ is Lipschitz continuous with Lipschitz constant bounded by $1 / \lambda$.

### 3.1. Regularized equation

Let us introduce the family of operators indexed by $\lambda>0$

$$
\begin{aligned}
B_{\lambda}: \Omega \times[0, T] \times H & \longrightarrow \mathscr{L}^{2}(U ; H) \\
(\omega, t, x) & \longmapsto B\left(\omega, t, J_{\lambda}(\omega, t, x)\right) .
\end{aligned}
$$

Since $A$, hence also $J_{\lambda}$, and $B$ are progressively measurable, $B_{\lambda}$ is progressively measurable as well for every $\lambda>0$. Moreover, $B_{\lambda}$ is Lipschitz continuous in its third argument, uniformly with respect to the other ones: in fact, thanks to assumption (III), one has

$$
\begin{aligned}
\frac{1}{2} \| B_{\lambda}(\omega, & t, x)-B_{\lambda}(\omega, t, y) \|_{\mathscr{L}^{2}(U ; H)}^{2} \\
= & \frac{1}{2}\left\|B\left(\omega, t, J_{\lambda}(\omega, t, x)\right)-B\left(\omega, t, J_{\lambda}(\omega, t, y)\right)\right\|_{\mathscr{L}^{2}(U ; H)}^{2} \\
\leqslant & \left\langle A\left(\omega, t, J_{\lambda}(\omega, t, x)\right)-A\left(\omega, t, J_{\lambda}(\omega, t, y)\right), J_{\lambda}(\omega, t, x)-J_{\lambda}(\omega, t, y)\right\rangle \\
& \quad+c_{2}\left\|J_{\lambda}(\omega, t, x)-J_{\lambda}(\omega, t, y)\right\|^{2} \\
= & \left\langle\tilde{A}_{\lambda}(\omega, t, x)-\tilde{A}_{\lambda}(\omega, t, y), x-y\right\rangle \leqslant \frac{1}{\lambda}\|x-y\|^{2}
\end{aligned}
$$

for every $\omega \in \Omega, t \in[0, T]$, and $x, y \in H$.
Therefore, since $\tilde{A}_{\lambda}$ is also a Lipschitz continuous operator on $H$, wellposedness results for stochastic differential equations on Hilbert spaces (see, e.g., $[10, \S 34]$ ) yields the existence (and uniqueness) of a predictable process $X_{\lambda} \in L^{2}(\Omega ; C([0, T] ; H))$ such that
(3.1) $X_{\lambda}+\int_{0}^{\cdot} \tilde{A}_{\lambda}\left(s, X_{\lambda}(s)\right) d s=X_{0}+c_{2} \int_{0} X_{\lambda}(s) d s+\int_{0} B_{\lambda}\left(s, X_{\lambda}(s)\right) d W(s)$.

### 3.2. A priori estimates

We are going to prove estimates on the family of solutions $\left(X_{\lambda}\right)$ to the regularized equations obtained above that are uniform with respect to $\lambda$.

The integration by parts formula for Hilbert space valued semimartingales yields

$$
\begin{aligned}
\frac{1}{2}\left\|X_{\lambda}\right\|^{2} & +\int_{0}\left\langle\tilde{A}_{\lambda}\left(s, X_{\lambda}(s)\right), X_{\lambda}(s)\right\rangle d s-\frac{1}{2} \int_{0}\left\|B_{\lambda}\left(s, X_{\lambda}(s)\right)\right\|_{\mathscr{L}^{2}(U ; H)}^{2} d s \\
& =\frac{1}{2}\left\|X_{0}\right\|^{2}+c_{2} \int_{0}\left\|X_{\lambda}(s)\right\|^{2} d s+\int_{0}^{\cdot} X_{\lambda}(s) B_{\lambda}\left(s, X_{\lambda}(s)\right) d W(s)
\end{aligned}
$$

where, in the stochastic integral, $X_{\lambda}$ is treated as a process taking values in the dual of $H$.

From now on we shall occasionally suppress the explicit indication of the dependence on $\omega$ and $t$ for processes and operators for notational simplicity. Note that, by definition of Yosida approximation,

$$
\begin{align*}
\left\langle\tilde{A}_{\lambda} x, x\right\rangle=\left\langle\tilde{A} J_{\lambda} x, x\right\rangle & =\left\langle\tilde{A} J_{\lambda} x, J_{\lambda} x\right\rangle+\left\langle\tilde{A} J_{\lambda} x, x-J_{\lambda} x\right\rangle  \tag{3.2}\\
& =\left\langle\tilde{A} J_{\lambda} x, J_{\lambda} x\right\rangle+\lambda\left\|\tilde{A}_{\lambda} x\right\|^{2}
\end{align*}
$$

for every $x \in H$. Therefore, denoting the norm of $\mathscr{L}^{2}(U ; H)$ by $\|\cdot\|_{2}$ for brevity, one has, thanks to assumption (III),

$$
\begin{aligned}
\left\langle\tilde{A}_{\lambda} x, x\right\rangle-\frac{1}{2}\left\|B_{\lambda}(x)\right\|_{2}^{2} & =\left\langle\tilde{A} J_{\lambda} x, J_{\lambda} x\right\rangle-\frac{1}{2}\left\|B\left(J_{\lambda} x\right)\right\|_{2}^{2}+\lambda\left\|\tilde{A}_{\lambda} x\right\|^{2} \\
& \geqslant c_{1}\left\|J_{\lambda} x\right\|_{V}^{p}+\lambda\left\|\tilde{A}_{\lambda} x\right\|^{2}-f
\end{aligned}
$$

which in turn yields

$$
\begin{gathered}
\frac{1}{2}\left\|X_{\lambda}\right\|^{2}+c_{1} \int_{0}\left\|J_{\lambda} X_{\lambda}(s)\right\|_{V}^{p} d s+\lambda \int_{0}\left\|\tilde{A}_{\lambda} X_{\lambda}(s)\right\|^{2} d s \\
\leqslant \\
\frac{1}{2}\left\|X_{0}\right\|^{2}+\int_{0} f(s) d s+c_{2} \int_{0}\left\|X_{\lambda}(s)\right\|^{2} d s \\
+\int_{0} X_{\lambda}(s) B_{\lambda}\left(X_{\lambda}(s)\right) d W(s)
\end{gathered}
$$

where the stochastic integral on the right-hand side is a martingale because $B_{\lambda}$ is Lipschitz continuous. In particular, taking expectation on both sides,

$$
\begin{gather*}
\frac{1}{2} \mathbb{E}\left\|X_{\lambda}(t)\right\|^{2}+c_{1} \mathbb{E} \int_{0}^{t}\left\|J_{\lambda} X_{\lambda}(s)\right\|_{V}^{p} d s+\lambda \mathbb{E} \int_{0}^{t}\left\|\tilde{A}_{\lambda} X_{\lambda}(s)\right\|^{2} d s \\
\leqslant \frac{1}{2} \mathbb{E}\left\|X_{0}\right\|^{2}+\mathbb{E} \int_{0}^{t} f(s) d s+c_{2} \mathbb{E} \int_{0}^{t}\left\|X_{\lambda}(s)\right\|^{2} d s \tag{3.3}
\end{gather*}
$$

for all $t \in[0, T]$. This implies that, for any interval $\left[0, T_{0}\right] \subseteq[0, T]$,

$$
\sup _{t \leqslant T_{0}} \mathbb{E}\left\|X_{\lambda}(t)\right\|^{2} \leqslant \mathbb{E}\left\|X_{0}\right\|^{2}+2 \mathbb{E} \int_{0}^{T} f(s) d s+2 c_{2} \mathbb{E} \int_{0}^{T_{0}}\left(\sup _{r \leqslant s} \mathbb{E}\left\|X_{\lambda}(r)\right\|^{2}\right) d s
$$

so that, by Gronwall's inequality, $\left(X_{\lambda}\right)$ is bounded in $C\left([0, T] ; L^{2}(\Omega ; H)\right)$. From this and (3.3) it immediately follows that $\left(\lambda^{1 / 2} \tilde{A}_{\lambda} X_{\lambda}\right)$ is bounded in $L^{2}\left(\Omega ; L^{2}(0, T ; H)\right)$ and that $\left(J_{\lambda} X_{\lambda}\right)$ is bounded in $L^{p}\left(\Omega ; L_{\tilde{A}}^{p}(0, T ; V)\right)$.
The latter implies, thanks to assumption (IV), that $\left(\tilde{A}_{\lambda} X_{\lambda}\right)$ is bounded in $L^{q}\left(\Omega ; L^{q}\left(0, T ; V^{\prime}\right)\right)$. Moreover, since, by assumption (III),

$$
\frac{1}{2}\left\|B_{\lambda}(x)\right\|_{2}^{2}=\frac{1}{2}\left\|B\left(J_{\lambda} x\right)\right\|_{2}^{2} \leqslant\left\langle\tilde{A} J_{\lambda} x, J_{\lambda} x\right\rangle+f=\left\langle\tilde{A}_{\lambda} x, J_{\lambda} x\right\rangle+f
$$

Hölder's inequality yields

$$
\begin{array}{r}
\mathbb{E} \int_{0}^{T}\left\|B_{\lambda}\left(X_{\lambda}(s)\right)\right\|_{2}^{2} d s \leqslant \mathbb{E} \int_{0}^{T}\left\|\tilde{A}_{\lambda} X_{\lambda}\right\|_{V^{\prime}}\left\|J_{\lambda} X_{\lambda}\right\|_{V}+\|f\|_{L^{1}(\Omega \times[0, T])} \\
\leqslant\left\|\tilde{A}_{\lambda} X_{\lambda}\right\|_{L^{q}\left(\Omega ; L^{q}\left(0, T ; V^{\prime}\right)\right)}\left\|J_{\lambda} X_{\lambda}\right\|_{L^{p}\left(\Omega ; L^{p}(0, T ; V)\right)} \\
+\|f\|_{L^{1}(\Omega \times[0, T])}
\end{array}
$$

thus proving that $\left(B_{\lambda}\left(X_{\lambda}\right)\right)$ is bounded in $L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right)$.

### 3.3. Construction of a solution and its uniqueness

The boundedness of various families of processes indexed by $\lambda$ obtained above imply, by well-known compactness properties in weak and weak* topologies, that there exist measurable and adapted processes

$$
\begin{aligned}
& X \in L^{\infty}\left(0, T ; L^{2}(\Omega ; H)\right) \\
& \bar{X} \in L^{p}\left(\Omega ; L^{p}(0, T ; V)\right) \\
& Y \in L^{q}\left(\Omega ; L^{q}\left(0, T ; V^{\prime}\right)\right) \\
& G \in L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right)
\end{aligned}
$$

such that

$$
\begin{aligned}
X_{\lambda} \longrightarrow X & \text { weakly* in } L^{\infty}\left(0, T ; L^{2}(\Omega ; H)\right), \\
J_{\lambda} X_{\lambda} \longrightarrow \bar{X} & \text { weakly in } L^{p}\left(\Omega ; L^{p}(0, T ; V)\right), \\
\lambda \tilde{A}_{\lambda} X_{\lambda} \longrightarrow 0 & \text { in } L^{2}\left(\Omega ; L^{2}(0, T ; H)\right), \\
\tilde{A}_{\lambda} X_{\lambda} \longrightarrow Y & \text { weakly in } L^{q}\left(\Omega ; L^{q}\left(0, T ; V^{\prime}\right)\right), \\
B_{\lambda}\left(\cdot, X_{\lambda}\right) \longrightarrow G & \text { weakly in } L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right) .
\end{aligned}
$$

Since, by definition of Yosida approximation, $X_{\lambda}-J_{\lambda} X_{\lambda}=\lambda \tilde{A}_{\lambda} X_{\lambda}$, the boundedness of $\left(\lambda^{1 / 2} \tilde{A}_{\lambda} X_{\lambda}\right)$ in $L^{2}\left(\Omega ; L^{2}(0, T ; H)\right)$ implies that

$$
X_{\lambda}-J_{\lambda} X_{\lambda} \longrightarrow 0 \quad \text { in } L^{2}\left(\Omega ; L^{2}(0, T ; H)\right)
$$

hence $X$ and $\bar{X}$ are equal $\mathbb{P} \otimes$ Leb-a.e. in $\Omega \times[0, T]$ and belong to

$$
L^{\infty}\left(0, T ; L^{2}(\Omega ; H)\right) \cap L^{p}\left(\Omega ; L^{p}(0, T ; V)\right)
$$

Recalling that the linear operator

$$
\begin{aligned}
L^{r}\left(\Omega ; L^{1}(0, T ; E)\right) & \longrightarrow L^{r}(\Omega ; C([0, T] ; E)) \\
w & \longmapsto \int_{0} w(s) d s
\end{aligned}
$$

is continuous for any $r \in[1, \infty[$ and any Banach space $E$, hence also weakly continuous, it follows from the above convergence results that

$$
\begin{aligned}
\int_{0}^{\cdot} X_{\lambda}(s) d s & \longrightarrow \int_{0} X(s) d s \\
\int_{0} \tilde{A}_{\lambda} X_{\lambda}(s) d s & \text { weakly in } L^{p}(\Omega ; C([0, T] ; V)) \\
\int_{0} Y(s) d s & \text { weakly in } L^{q}\left(\Omega ; C\left([0, T] ; V^{\prime}\right)\right)
\end{aligned}
$$

Similarly, since the stochastic integral operator

$$
\begin{aligned}
L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right. & \longrightarrow L^{2}(\Omega ; C([0, T] ; H)) \\
C & \longmapsto \int_{0} C(s) d W(s)
\end{aligned}
$$

is linear and continuous, hence weakly continuous, one gets that

$$
\int_{0} B_{\lambda}\left(s, X_{\lambda}(s)\right) d W(s) \longrightarrow \int_{0} G(s) d W(s) \quad \text { weakly in } L^{2}(\Omega ; C([0, T] ; H))
$$

This implies, taking the weak limit in the space $L^{q}\left(\Omega ; C\left([0, T] ; V^{\prime}\right)\right)$ as $\lambda \rightarrow 0$ in equation (3.1),

$$
\begin{equation*}
X+\int_{0} Y(s) d s=X_{0}+c_{2} \int_{0} X(s) d s+\int_{0} G(s) d W(s) \quad \text { in } V^{\prime} \tag{3.4}
\end{equation*}
$$

We are going to show, using an adaptation of a classical argument from the theory of maximal monotone operators (cf. [5] as well as [7, §4.2]), that $Y=$ $A(\cdot, X)+c_{2} X$ and $G=B(X)$. Itô's formula for the square of the $H$-norm applied to (3.1) yields

$$
\begin{aligned}
& \left.\frac{1}{2}\left\|e^{-c_{2} \cdot} X_{\lambda}\right\|^{2}+\int_{0} e^{-2 c_{2} s}\left\langle\tilde{A}_{\lambda} X_{\lambda}(s)\right), X_{\lambda}(s)\right\rangle d s \\
& \quad-\frac{1}{2} \int_{0} e^{-c_{2} s}\left\|B_{\lambda}\left(X_{\lambda}(s)\right)\right\|_{\mathscr{L}^{2}(U ; H)}^{2} d s \\
& \quad=\frac{1}{2}\left\|X_{0}\right\|^{2}+\int_{0} e^{-c_{2} s} X_{\lambda}(s) B_{\lambda}\left(X_{\lambda}(s)\right) d W(s)
\end{aligned}
$$

Since the stochastic integral on the right-hand side is a martingale, as seen above, one has, applying inequality (3.2) and taking expectations on both sides,

$$
\begin{aligned}
& \frac{1}{2} \mathbb{E}\left\|e^{-c_{2} \cdot} X_{\lambda}\right\|^{2}-\frac{1}{2} \mathbb{E}\left\|X_{0}\right\|^{2} \\
& \quad \leqslant \mathbb{E} \int_{0} e^{-2 c_{2} s}\left(-\left\langle\tilde{A} J_{\lambda} X_{\lambda}(s), J_{\lambda} X_{\lambda}(s)\right\rangle+\frac{1}{2}\left\|B\left(J_{\lambda} X_{\lambda}(s)\right)\right\|_{\mathscr{L}^{2}(U ; H)}^{2}\right) d s
\end{aligned}
$$

Let $\varphi \in L^{2}(\Omega ; C([0, T] ; H)) \cap L^{p}\left(\Omega ; L^{p}(0, T ; V)\right)$.
Thanks to assumptions (III)-(IV), one has that $\tilde{A} \varphi \in L^{q}\left(\Omega ; L^{q}\left(0, T ; V^{\prime}\right)\right)$ and $B(\varphi) \in L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right)$, so that the term within parentheses on right-hand side can be rewritten as

$$
\begin{aligned}
-\left\langle\tilde{A} J_{\lambda}\right. & \left.X_{\lambda}-\tilde{A} \varphi, J_{\lambda} X_{\lambda}-\varphi\right\rangle+\frac{1}{2}\left\|B\left(J_{\lambda} X_{\lambda}\right)-B(\varphi)\right\|_{\mathscr{L}^{2}(U ; H)}^{2} \\
& -\left\langle\tilde{A} \varphi, J_{\lambda} X_{\lambda}\right\rangle-\left\langle\tilde{A} J_{\lambda} X_{\lambda}-\tilde{A} \varphi, \varphi\right\rangle \\
& -\frac{1}{2}\|B(\varphi)\|_{\mathscr{L}^{2}(U ; H)}^{2}+\left\langle B\left(J_{\lambda} X_{\lambda}\right), B(\varphi)\right\rangle_{\mathscr{L}^{2}(U ; H)}
\end{aligned}
$$

where the sum of the two terms in the first row is negative by assumption (III). One is thus left with

$$
\begin{aligned}
& \frac{1}{2} \mathbb{E}\left\|e^{-c_{2} \cdot} X_{\lambda}\right\|^{2}-\frac{1}{2} \mathbb{E}\left\|X_{0}\right\|^{2} \\
& \leqslant \mathbb{E} \int_{0} e^{-2 c_{2} s}\left(-\left\langle\tilde{A} \varphi(s), J_{\lambda} X_{\lambda}(s)\right\rangle-\left\langle\tilde{A} J_{\lambda} X_{\lambda}(s)-\tilde{A} \varphi(s), \varphi(s)\right\rangle\right. \\
& \left.\quad-\frac{1}{2}\|B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}+\left\langle B\left(J_{\lambda} X_{\lambda}(s)\right), B(\varphi(s))\right\rangle_{\mathscr{L}^{2}(U ; H)}\right) d s
\end{aligned}
$$

from which it follows that, for every nonnegative $\psi \in L^{\infty}(0, T)$,

$$
\begin{aligned}
& \frac{1}{2} \mathbb{E} \int_{0}^{T} \psi(t)\left(e^{-2 c_{2} t}\left\|X_{\lambda}(t)\right\|^{2}-\left\|X_{0}\right\|^{2}\right) d t \\
& \leqslant \mathbb{E} \int_{0}^{T} \psi(t)\left(\int _ { 0 } ^ { t } e ^ { - 2 c _ { 2 } s } \left(-\left\langle\tilde{A} \varphi(s), J_{\lambda} X_{\lambda}(s)\right\rangle-\left\langle\tilde{A} J_{\lambda} X_{\lambda}(s)-\tilde{A} \varphi(s), \varphi(s)\right\rangle\right.\right. \\
& \left.\left.\quad-\frac{1}{2}\|B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}+\left\langle B\left(J_{\lambda} X_{\lambda}(s)\right), B(\varphi(s))\right\rangle_{\mathscr{L}^{2}(U ; H)}\right) d s\right) d t
\end{aligned}
$$

By the weak lower semicontinuity of the norm in $L^{2}\left(\Omega ; L^{2}(0, T ; H)\right)$, one has

$$
\begin{aligned}
\mathbb{E} \int_{0}^{T} \psi(t)\left(e^{-2 c_{2} t}\|X(t)\|^{2}\right. & \left.-\left\|X_{0}\right\|^{2}\right) d t \\
& \leqslant \liminf _{\lambda \rightarrow 0} \mathbb{E} \int_{0}^{T} \psi(t)\left(e^{-2 c_{2} t}\left\|X_{\lambda}(t)\right\|^{2}-\left\|X_{0}\right\|^{2}\right) d t
\end{aligned}
$$

Moreover, since $B(\varphi) \in L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right)$, the weak convergence results proved above yield

$$
\begin{aligned}
& \lim _{\lambda \rightarrow 0} \mathbb{E} \int_{0}^{T} \psi(t)\left(\int _ { 0 } ^ { t } e ^ { - 2 c _ { 2 } s } \left(-\left\langle\tilde{A} \varphi(s), J_{\lambda} X_{\lambda}(s)\right\rangle-\left\langle\tilde{A} J_{\lambda} X_{\lambda}(s)-\tilde{A} \varphi(s), \varphi(s)\right\rangle\right.\right. \\
&\left.\left.-\frac{1}{2}\|B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}+\left\langle B\left(J_{\lambda} X_{\lambda}(s)\right), B(\varphi(s))\right\rangle_{\mathscr{L}^{2}(U ; H)}\right) d s\right) d t \\
&=\mathbb{E} \int_{0}^{T} \psi(t)\left(\int_{0}^{t} e^{-2 c_{2} s}(-\langle\tilde{A}(\varphi(s)), X(s)\rangle-\langle Y(s)-\tilde{A}(\varphi(s)), \varphi(s)\rangle\right. \\
&\left.\left.\quad-\frac{1}{2}\|B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}+\langle G(s), B(\varphi(s))\rangle_{\mathscr{L}^{2}(U ; H)}\right) d s\right) d t
\end{aligned}
$$

We then deduce that

$$
\begin{aligned}
& \frac{1}{2} \mathbb{E} \int_{0}^{T} \psi(t)\left(e^{-2 c_{2} t}\|X(t)\|^{2}-\left\|X_{0}\right\|^{2}\right) d t \\
& \quad \leqslant \mathbb{E} \int_{0}^{T} \psi(t)\left(\int_{0}^{t} e^{-2 c_{2} s}(-\langle\tilde{A} \varphi(s), X(s)\rangle-\langle Y(s)-\tilde{A} \varphi(s), \varphi(s)\rangle\right.
\end{aligned}
$$

$$
\left.\left.-\frac{1}{2}\|B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}+\langle G(s), B(\varphi(s))\rangle_{\mathscr{L}^{2}(U ; H)}\right) d s\right) d t
$$

Itô's formula (see [5] or $[7, \S 4]$ ) applied to the limit equation (3.4) implies that there exists a modification of $X$, denoted by the same symbol for simplicity, such that $X \in L^{2}(\Omega ; C([0, T] ; H))$ and

$$
\begin{gathered}
\frac{1}{2} e^{-2 c_{2} \cdot}\|X\|^{2}+\int_{0} e^{-2 c_{2} s}\langle Y(s), X(s)\rangle d s-\frac{1}{2} \int_{0} e^{-2 c_{2} s}\|G(s)\|_{\mathscr{L}^{2}(U ; H)}^{2} d s \\
=\frac{1}{2}\left\|X_{0}\right\|^{2}+\int_{0} e^{-2 c_{2} s} X(s) G(s) d W(s)
\end{gathered}
$$

Substituting this identity on the left-hand side of the last inequality and rearranging the terms yields

$$
\begin{aligned}
\mathbb{E} \int_{0}^{T} \psi(t) \int_{0}^{t} e^{-2 c_{2} s} & (\langle Y(s)-\tilde{A} \varphi(s), \varphi(s)-X(s)\rangle \\
& \left.+\frac{1}{2}\|G(s)-B(\varphi(s))\|_{\mathscr{L}^{2}(U ; H)}^{2}\right) d s d t \leqslant 0
\end{aligned}
$$

Choosing $\varphi=X$ immediately yields $G=B(X)$, hence also, in particular,

$$
\mathbb{E} \int_{0}^{T} \psi(t) \int_{0}^{t} e^{-2 c_{2} s}\langle Y(s)-\tilde{A} \varphi(s), \varphi(s)-X(s)\rangle d s d t \leqslant 0
$$

Let $\delta \in \mathbb{R}_{+}, v \in V$, and $\bar{\varphi} \in L^{\infty}(\Omega \times[0, T])$. Choosing $\varphi=X+\delta \bar{\varphi} v$ and taking the limit as $\delta \rightarrow 0$ yields

$$
\mathbb{E} \int_{0}^{T} \psi(t) \int_{0}^{t} e^{-2 c_{2} s} \bar{\varphi}(s)\langle Y(s)-\tilde{A} \varphi(s), v\rangle d s d t \leqslant 0
$$

By a classical localisation argument, recalling that $\psi$ and $\bar{\varphi}$ have been chosen arbitrarily, one has

$$
\langle Y-\tilde{A} \varphi, v\rangle \leqslant 0 \quad \forall v \in V
$$

a.e. in $\Omega \times[0, T]$. Since $v$ has also been chosen arbitrarily in $V$, it follows that

$$
\langle Y-\tilde{A} \varphi, X-\varphi\rangle \geqslant 0
$$

a.e. in $\Omega \times[0, T]$. Since $\tilde{A}$ is maximal monotone, this implies that $Y=\tilde{A} X$ a.e. in $\Omega \times[0, T]$. We have thus shown that $X$ is a strong solution to equation (1.1).

In order to prove the Lipschitz continuous dependence of the solution with respect to the initial datum (from which uniqueness follows), let $X_{1}, X_{2}$ be strong solutions to (1.1), in the sense of Theorem 2.2 , with initial data $X_{0}^{1}$ and $X_{0}^{2}$, respectively. Itô formula (as in [5] or [7]) yields

$$
\begin{aligned}
& \frac{1}{2}\left\|\left(X_{1}-X_{2}\right)\right\|_{H}^{2} \\
& \quad+\int_{0}\left(\left\langle A X_{1}-A X_{2}, X_{1}-X_{2}\right\rangle-\frac{1}{2}\left\|B\left(X_{1}\right)-B\left(X_{2}\right)\right\|_{\mathscr{L}^{2}(U ; H)}^{2}\right)(s) d s \\
& \quad=\frac{1}{2}\left\|X_{0}^{1}-X_{0}^{2}\right\|^{2}+\int_{0}\left(\left(X_{1}-X_{2}\right)\left(B\left(X_{1}\right)-B\left(X_{2}\right)\right)\right)(s) d W(s),
\end{aligned}
$$

where the stochastic integral on the right-hand side is a martingale because $X_{1}, X_{2} \in L^{2}(\Omega ; C([0, T] ; H))$ and $B\left(X_{1}\right), B\left(X_{2}\right) \in L^{2}\left(\Omega ; L^{2}\left(0, T ; \mathscr{L}^{2}(U ; H)\right)\right)$. Therefore, taking expectations on both sides and using assumption (III), $\mathbb{E}\left\|\left(X_{1}-X_{2}\right)(t)\right\|^{2} \leqslant \mathbb{E}\left\|X_{0}^{1}-X_{0}^{2}\right\|^{2}+2 c_{2} \int_{0}^{t} \mathbb{E}\left\|\left(X_{1}-X_{2}\right)(s)\right\|^{2} d s \quad \forall t \in[0, T]$,
from which the conclusion follows thanks to Gronwall's inequality. The proof of Theorem 2.2 is thus complete.

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