# The stochastic $p(\omega, t, x)$-Laplace equation with cylindrical Wiener process 

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We propose the analysis of a non-linear parabolic problem of $p(\omega, t, x)$-Laplace type in the framework of Orlicz Lebesgue and Sobolev spaces with variable random exponents and a stochastic forcing by a cylindrical Wiener process. We give a result of well-posedness: existence, uniqueness and stability of the solution, for additive and multiplicative problems.

## 1. Introduction

We are interested in a result of existence and uniqueness of the solution to the problem:

$$
(P, H) \begin{cases}d u-\Delta_{p(\cdot)} u d t=H(\cdot, u) d W & \text { in } \Omega \times(0, T) \times D \\ u=0 & \text { on } \Omega \times(0, T) \times \partial D \\ u(0, \cdot)=u_{0} & \text { in } L^{2}(D)\end{cases}
$$

where $T>0, D \subset \mathbb{R}^{d}$ is a bounded Lipschitz domain, $Q:=(0, T) \times D$ and $(\Omega, \mathcal{F}, P)$ is a classical Wiener space endowed with a normal filtration $\left(\mathcal{F}_{t}\right)_{t \geq 0}$.

$$
H: \Omega \times(0, T) \times \mathbb{R} \rightarrow H S\left(L^{2}(D)\right), \quad(\omega, t, \lambda) \mapsto H(\omega, t, \lambda)
$$

is a Carathéodory function, continuous with respect to $\lambda$, progressively measurable with respect to $\mathcal{F}_{t}$ and square integrable with respect to $d P \otimes d t$, where $H S\left(L^{2}(D)\right)$ denotes the space of Hilbert-Schmidt operators on $L^{2}(D)$ with values in $L^{2}(D)$. We will give the precise assumptions on $H$ below. $\Delta_{p(\cdot)} u=$

[^0]$\operatorname{div}\left(|\nabla u|^{p(\omega, t, x)-2} \nabla u\right)$ denotes the $p$-Laplace operator with a variable exponent $p: \Omega \times Q \rightarrow(1, \infty)$ satisfying the following conditions:
( $p 1$ ) $1<p^{-}:=\operatorname{ess}_{\inf }^{(\omega, t, x)}$ $p(\omega, t, x) \leq p^{+}:=\operatorname{esssup}_{(\omega, t, x)} p(\omega, t, x)<\infty$,
( $p 2$ ) $\omega$ a.s. in $\Omega,(t, x) \mapsto p(\omega, t, x)$, is log-Hölder continuous, i.e. there exists $C \geq 0$ (which may depend on $\omega$ ) such that, for all $(t, x),(s, y) \in Q$,
\[

$$
\begin{equation*}
|p(\omega, t, x)-p(\omega, s, y)| \leq \frac{C}{\ln \left(e+\frac{1}{|(t, x)-(s, y)|}\right)} \tag{1}
\end{equation*}
$$

\]

( $p 3$ ) progressive measurability of the variable exponent, i.e.

$$
\Omega \times[0, t] \times D \ni(\omega, s, x) \mapsto p(\omega, s, x)
$$

is $\mathcal{F}_{t} \times \mathcal{B}(0, t) \times \mathcal{B}(D)$-measurable for all $0 \leq t \leq T$.
For an orthonormal basis $\left(e_{k}\right)$ of $L^{2}(D)$ and $\left(\beta_{k}(t)\right)$ a family of independent, real-valued Brownian motions adapted to $\left(\mathcal{F}_{t}\right)$, we (formally) define the cylindrical Wiener process,

$$
\begin{equation*}
W(t):=\sum_{k=1}^{\infty} e_{k} \beta_{k}(t) \tag{2}
\end{equation*}
$$

It is well-known that the sum on the right-hand side of (2) does not converge in $L^{2}(D)$, therefore we have to give a meaning to (2) following the ideas of [2] and [7]: For $u=\sum_{k=1}^{\infty} u_{k} e_{k}$ and $v=\sum_{k=1}^{\infty} v_{k} e_{k}$

$$
(u, v)_{U}:=\sum_{k=1}^{\infty} \frac{u_{k} v_{k}}{k^{2}}
$$

is a scalar product on $L^{2}(D)$. Now we define the (bigger) Hilbert space $U$ as the completion of $L^{2}(D)$ with respect to the norm $\|\cdot\|_{U}$ induced by $(\cdot, \cdot)_{U}$. It is then easy to see that $\left(k e_{k}\right)$ is an orthonormal basis of $U$. Note that

$$
\begin{equation*}
W(t)=\sum_{k=1}^{\infty} e_{k} \beta_{k}(t)=\sum_{k=1}^{\infty} \frac{1}{k} k e_{k} \beta_{k}(t) \tag{3}
\end{equation*}
$$

and therefore $W(t)$ can be interpreted as a $Q$-Wiener process with covariance Matrix $Q=\operatorname{diag}\left(\frac{1}{k^{2}}\right)$ and values in $U$. Since $Q^{\frac{1}{2}}(U)=L^{2}(D)$, for all square integrable and predictable $\Phi: \Omega \times(0, T) \rightarrow H S\left(L^{2}(D)\right)$ the stochastic integral with respect to the cylindrical Wiener process $W(t)$ can be defined by

$$
\int_{0}^{t} \Phi d W=\sum_{k=1}^{\infty} \int_{0}^{t} \frac{1}{k} \Phi\left(k e_{k}\right) d \beta_{k}=\sum_{k=1}^{\infty} \int_{0}^{t} \Phi\left(e_{k}\right) d \beta_{k}
$$

In particular, $\Phi\left(e_{k}\right) \in N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ for all $k \in \mathbb{N}^{*}$.
Assume that $H: \Omega \times(0, T) \times L^{2}(D) \rightarrow H S\left(L^{2}(D)\right)$ is defined by

$$
H(\omega, t, u)\left(e_{k}\right)=\left\{x \mapsto h_{k}(\omega, t, x, u(x))\right\}
$$

where, for any $k \in \mathbb{N}^{*}, h_{k}: \Omega \times(0, T) \times \mathbb{R}^{d+1} \rightarrow \mathbb{R}$ is a Carathéodory function such that for all $\lambda \in \mathbb{R}$, $h_{k}(\cdot, \lambda) \in N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ and $\lambda \mapsto h_{k}(\omega, t, \lambda)$ is continuous $P_{T} \otimes \mathcal{L}^{d}$-a.e. where $P_{T}$ denotes the predictable $\sigma$-algebra and $\mathcal{L}$ the Lebesgue measure (see [2] for example). Moreover, for technical reasons,
(H1) There exist $C_{1}, C_{2} \geq 0$ and $C_{3} \in L^{1}(D)$ such that a.e. in $(\omega, t, x)$,

$$
\sum_{k=1}^{\infty}\left|h_{k}(\cdot, \lambda)\right|^{2} \leq C_{1}|\lambda|^{2}+C_{3}, \sum_{k=1}^{\infty}\left|h_{k}(\cdot, \lambda)-h_{k}(\cdot, \mu)\right|^{2} \leq C_{2}|\lambda-\mu|^{2} .
$$

In particular for any $u \in L^{2}(D)$ and for a.e. $(\omega, t) \in \Omega \times(0, T)$ thanks to (H1) we have

$$
\begin{aligned}
\|H(\omega, t, u)\|_{H S\left(L^{2}(D)\right)}^{2} & =\sum_{k=1}^{\infty}\left\|H(\omega, t, u)\left(e_{k}\right)\right\|_{L^{2}(D)}^{2} \\
& =\int_{D} \sum_{k=1}^{\infty}\left|h_{k}(\omega, t, x, u(x))\right|^{2} d x \leq\left\|C_{3}\right\|_{L^{1}(D)}+C_{1}\|u\|_{L^{2}(D)}^{2}
\end{aligned}
$$

and therefore $H(\omega, t, u)$ is a Hilbert-Schmidt operator.
Our aim in this paper is to extend the previous result published in [1] to the case of a random variable exponent and to a more general noise, here a cylindrical Wiener process. At the beginning, the methodology is close to the one presented in [1], then it has to be adapted to the new situation. The result is first proved in the additive case for a finite-dimensional Wiener process: as in [1], one considers a singular perturbation of the $p(\cdot)$-Laplace operator by a $q$-Laplace one ( $q$ being a big enough constant) with very regular additive integrands $H$ before passing to the limits on the perturbation, then on the regularization of the integrands. The result is then proved in the additive case for a general infinite-dimensional Wiener process, then in the multiplicative case by using a fixed point argument.

The organization of the paper is the following one: the next section presents the functional framework and the one after introduces the main result. The last section is dedicated to the proof of the main result.

## 2. Function spaces with variable exponent

The following function space serves as the variable exponent version of the classical Bochner space setting: there exists a full-measure set $\tilde{\Omega} \subset \Omega$ such that we can define

$$
X_{\omega}(Q):=\left\{u \in L^{2}(Q) \cap L^{1}\left(0, T ; W_{0}^{1,1}(D)\right) \mid \nabla u \in\left(L^{p(\omega,)}(Q)\right)^{d}\right\}
$$

which is a separable, reflexive Banach space for all $\omega \in \tilde{\Omega}$ with respect to the norm

$$
\|u\|_{X_{\omega}(Q)}=\|u\|_{L^{2}(Q)}+\|\nabla u\|_{L^{p(\omega,)}(Q)}
$$

$X_{\omega}(Q)$ is a parametrization by $\omega$ of the space

$$
X(Q):=\left\{u \in L^{2}(Q) \cap L^{1}\left(0, T ; W_{0}^{1,1}(D)\right) \mid \nabla u \in\left(L^{p(t, x)}(Q)\right)^{d}\right\}
$$

which has been introduced in [3] for the case of a variable exponent depending on $(t, x)$. For the basic properties of $X(Q)$, we refer to [3]. For $u \in X_{\omega}(Q)$, it follows directly from the definition that $u(t) \in$ $L^{2}(D) \cap W_{0}^{1,1}(D)$ for almost every $t \in(0, T)$. Moreover, from $\nabla u \in L^{p(\omega, \cdot)}(Q)$ and Fubini's theorem it follows that $\nabla u(t, \cdot)$ is in $L^{p(\omega, t, \cdot)}(D)$ a.e. in $(0, T)$.

Let us introduce the space

$$
\mathcal{E}:=\left\{u \in L^{2}(\Omega \times Q) \cap L^{p^{-}}\left(\Omega \times(0, T) ; W_{0}^{1, p^{-}}(D)\right) \mid \nabla u \in L^{p(\cdot)}(\Omega \times Q)\right\}
$$

which is a separable, reflexive Banach space with respect to the norm

$$
u \in \mathcal{E} \mapsto\|u\|_{\mathcal{E}}=\|u\|_{L^{2}(\Omega \times Q)}+\|\nabla u\|_{L^{p(\cdot)}(\Omega \times Q)} .
$$

Thanks to Fubini's theorem, $u \in \mathcal{E}$ implies that $u(\omega) \in X_{\omega}(Q)$ a.s. in $\Omega$ and, since Poincaré's inequality is available with respect to $(t, x)$, independently of $\omega, u \in \mathcal{E}$ implies also $u(\omega, t) \in L^{2}(D) \cap W_{0}^{1, p(\omega, t, \cdot)}(D)$ for almost all $(\omega, t) \in \Omega \times(0, T)$.

## 3. Main result

Definition 3.1. A solution to $(P, H)$ is a function $u \in L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right) \cap N_{w}^{2}\left(0, T ; L^{2}(D)\right) \cap \mathcal{E}$, such that, for almost every $\omega \in \Omega, u(0, \cdot)=u_{0}$, a.e. in $D$ and for all $t \in[0, T]$,

$$
u(t)-u_{0}-\int_{0}^{t} \Delta_{p(\cdot)} u d s=\int_{0}^{t} H(\cdot, u) d W,
$$

holds a.s. in $D$; or, equivalently, in the weak sense:

$$
\partial_{t}\left[u(t)-\int_{0}^{t} H(\cdot, u) d W\right]-\Delta_{p(\cdot)} u=0 \text { in } X_{\omega}^{\prime}(Q) .
$$

Theorem 3.1. There exists a unique solution to $(P, H)$. Moreover, if $u_{1}, u_{2}$ are the solutions to $\left(P, H_{1}\right)$, $\left(P, H_{2}\right)$ respectively, then:

$$
\begin{align*}
& E\left(\sup _{t \in[0, T]}\left\|\left(u_{1}-u_{2}\right)(t)\right\|_{L^{2}(D)}^{2}\right)+E \int_{Q}\left(\left|\nabla u_{1}\right|^{p(\cdot)-2} \nabla u_{1}-\left|\nabla u_{2}\right|^{p(\cdot)-2} \nabla u_{2}\right) \cdot \nabla\left(u_{1}-u_{2}\right) d(t, x) \\
\leq & C E \int_{0}^{T}\left\|H_{1}\left(\cdot, u_{1}\right)-H_{2}\left(\cdot, u_{2}\right)\right\|_{H S\left(L^{2}(D)\right)}^{2} d t . \tag{4}
\end{align*}
$$

## 4. Proof of the main result

Notation: for a square integrable and predictable process $\Phi: \Omega \times(0, T) \rightarrow H S\left(L^{2}(D)\right)$ and $N \in \mathbb{N}^{*}$, we define the predictable and square integrable process $\Phi_{N}: \Omega \times(0, T) \rightarrow H S\left(L^{2}(D)\right)$ by $\Phi_{N}\left(e_{k}\right)=\Phi\left(e_{k}\right)$ for $k \leq N$ and $\phi_{N}\left(e_{k}\right)=0$ for $k>N$. Consequently

$$
\int_{0}^{T} \Phi_{N} d W=\sum_{k=1}^{N} \int_{0}^{T} \Phi\left(e_{k}\right) d \beta_{k} .
$$

Let us remark that this corresponds to the case of the finite-dimensional Wiener process: $W_{N}(t)=$ $\sum_{k=1}^{N} \beta_{k}(t) e_{k}$.

### 4.1. The result for nice processes

Let $S_{w}^{2}\left(0, T ; H_{0}^{j}(D)\right)$ be the subset of simple, predictable processes with values in $H_{0}^{j}(D)$ for sufficiently large values of $j$. Note that $S_{w}^{2}\left(0, T ; H_{0}^{j}(D)\right)$ is densely imbedded into $N_{w}^{2}\left(0, T ; L^{2}(D)\right)$. We will first prove the result when $\Phi_{N}\left(e_{k}\right) \in S_{w}^{2}\left(0, T ; H_{0}^{j}(D)\right)$ for all $k=1, \ldots N$. We will call such $\Phi_{N}$ a nice process in the sequel.

Proposition 4.1. For $q \geq \max \left(2, p^{+}\right), 0<\varepsilon \leq 1, N \in \mathbb{N}^{*}$ and a nice process $\Phi_{N}$ there exists

$$
u^{\varepsilon} \in L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right) \cap N_{w}^{2}\left(0, T ; L^{2}(D)\right) \cap L^{q}\left(\Omega \times(0, T) ; W_{0}^{1, q}(D)\right)
$$

and a set $\tilde{\Omega} \subset \Omega$ of total probability 1 on which $u(0, \cdot)=u_{0}$ a.e. in $D$ and

$$
\begin{equation*}
u^{\varepsilon}(t)-u_{0}-\int_{0}^{t}\left[\varepsilon \Delta_{q} u^{\varepsilon}+\Delta_{p(\cdot)} u^{\varepsilon}\right] d s=\int_{0}^{t} \Phi_{N} d W \tag{5}
\end{equation*}
$$

in $W^{-1, q^{\prime}}(D)$ for all $t \in[0, T]$.
Proof. For $q \geq \max \left(2, p^{+}\right)$and $\varepsilon>0$, the operator

$$
A: \Omega \times(0, T) \times W_{0}^{1, q}(D) \rightarrow W^{-1, q^{\prime}}(D), \quad A(\omega, t, u)=-\varepsilon \Delta_{q} u-\Delta_{p(\omega, t, x)} u
$$

is monotone with respect to $u$ for a.e. $(\omega, t) \in \Omega \times(0, T)$ and progressively measurable, i.e. for every $t \in[0, T]$ the mapping

$$
A: \Omega \times(0, t) \times W_{0}^{1, q}(D) \rightarrow W^{-1, q^{\prime}}(D), \quad(\omega, s, u) \mapsto A(\omega, s, u)
$$

is $\mathcal{F}_{t} \times \mathcal{B}(0, t) \times \mathcal{B}\left(W_{0}^{1, q}(D)\right)$-measurable. In particular, $-A$ satisfies the hypotheses of [5, Theorem 2.1, p. 1253], therefore for any $\varepsilon>0$ there exists a continuous process with values in $L^{2}(D)$ solution to the problem (5). Then, [2, Prop. 3.17 p. 84] and [5, Theorem 2.3, p. 1254] yield $u^{\varepsilon} \in L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$.

Proposition 4.2. For any nice process $\Phi_{N}$, there exist a unique function $u \in \mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$ and a full-measure set $\tilde{\Omega} \in \mathcal{F}$ such that for all $\omega \in \tilde{\Omega}$ we have $u(0, \cdot)=u_{0}$ a.e. in $D$ and

$$
\begin{equation*}
u(t)-u_{0}-\int_{0}^{t} \Delta_{p(\cdot)} u d s=\int_{0}^{t} \Phi_{N} d W \tag{6}
\end{equation*}
$$

holds a.e. in $D$ for all $t \in[0, T]$. In particular $u$ is a solution to $\left(P, \Phi_{N}\right)$ in the sense of Definition 3.1.
Proof. For the first part of the proof, mainly based on deterministic arguments, we can repeat the arguments of [1]: If we set $v^{\varepsilon}:=u^{\varepsilon}-\int_{0}^{t} \Phi_{N} d W$, such that $v^{\varepsilon}(0)=u_{0}$, then $u^{\varepsilon}$ satisfies (5), iff there exists a full-measure set $\tilde{\Omega} \in \mathcal{F}$ such that

$$
\begin{equation*}
\partial_{t} v^{\varepsilon}-\varepsilon \Delta_{q}\left(v^{\varepsilon}+\int_{0}^{t} \Phi_{N} d W\right)-\Delta_{p(\cdot)}\left(v^{\varepsilon}+\int_{0}^{t} \Phi_{N} d W\right)=0 \tag{7}
\end{equation*}
$$

in $L^{q^{\prime}}\left(0, T ; W^{-1, q^{\prime}}(D)\right)$ for all $\omega \in \tilde{\Omega}$. Testing (7) with $v^{\varepsilon}$ to get a priori estimates, we can use classical (monotonicity) arguments to conclude that pointwise for every $\omega \in \tilde{\Omega}$ we have the following convergence results, passing to a (not relabeled) subsequence if necessary when $\epsilon$ tends to 0 :
1.) $v^{\varepsilon} \rightharpoonup v$ in $X_{\omega}(Q)$ and $L^{\infty}\left(0, T ; L^{2}(D)\right)$ weak-*,
2.) for any $t, v^{\varepsilon}(t) \rightarrow v(t)$ in $L^{2}(D)$,
3.) $\int_{Q}\left|\nabla v^{\varepsilon}-\nabla v\right|^{p(\omega, t, x)} d x d t \rightarrow 0$.

Then, passing to the limit in the singular perturbation, $v$ satisfies the problem

$$
\partial_{t} v-\Delta_{p(\cdot)}\left(v+\int_{0}^{t} \Phi_{N} d W\right)=0
$$

In particular, $\partial_{t} v \in X_{\omega}^{\prime}(Q)$ (see [3]) and $v \in W_{\omega}(Q)$ where one denotes by

$$
W_{\omega}(Q):=\left\{v \in X_{\omega}(Q) \mid \partial_{t} v \in X_{\omega}^{\prime}(Q)\right\} .
$$

Thanks to [3], $W_{\omega}(Q) \hookrightarrow C\left([0, T] ; L^{2}(D)\right)$ with a continuity constant depending only on $T$ and the timeintegration by parts formula is available. Thus, $v \in C\left([0, T] ; L^{2}(D)\right)$ and $v$ is a solution of the above problem in $W_{\omega}(Q)$, for the initial condition $u_{0}$. Since this solution is unique, no subsequence is needed in the above limits. Then, the above convergence yields for all $\omega \in \tilde{\Omega}$ :
1.) $u^{\varepsilon} \rightarrow u$ in $L^{2}\left(0, T ; L^{2}(D)\right)$ with $\partial_{t}\left[u-\int_{0}^{*} \Phi_{N} d W\right] \in X_{\omega}^{\prime}(Q)$,
2.) for any $t, u^{\varepsilon}(t) \rightarrow u(t)$ in $L^{2}(D)$,
3.) $\Delta_{p(\omega, t, x)} u^{\varepsilon} \rightharpoonup \Delta_{p(\omega, t, x)} u$ in $X_{\omega}^{\prime}(Q)$,
4.) $\int_{Q}\left|\nabla u^{\varepsilon}-\nabla u\right|^{p(\omega, t, x)} d x d t \rightarrow 0$.

We continue with the argumentation as in [1]: from the previous convergence results, the a priori estimates and since $\nabla \Phi_{N}$ is bounded, we get uniform estimates that allow us to use Lebesgue Dominated Convergence theorem and therefore it follows that

$$
\begin{equation*}
\forall t, u^{\varepsilon}(t) \rightarrow u(t) \text { in } L^{2}\left(\Omega ; L^{2}(D)\right) \quad \text { and } \quad u^{\varepsilon} \rightarrow u \text { in } \mathcal{E} \tag{8}
\end{equation*}
$$

Note that the above limits in $L^{2}\left(\Omega ; L^{2}(D)\right)$ and $L^{2}\left(\Omega ; L^{2}(Q)\right)$ are standard results obtained in classical Bochner spaces, but the measurability of $\nabla u$ with respect to $d(t, x) \otimes d P$ deserves our attention. Since $\nabla u^{\varepsilon}$ and $\nabla u^{\epsilon^{\prime}}$ are globally measurable functions, Lebesgue Dominated Convergence theorem, together with $a$ priori estimates yield

$$
E \int_{Q}\left|\nabla u^{\varepsilon}-\nabla u^{\varepsilon^{\prime}}\right|^{p(\omega, t, x)} d x d t \rightarrow 0
$$

and thus, $\left(\nabla u^{\varepsilon}\right)$ is a Cauchy sequence in $L^{p(\cdot)}(\Omega \times Q)$ and therefore a converging sequence. It is then a direct consequence to see that $\nabla u$ is the limit in $L^{p(\cdot)}(\Omega \times Q)$ of $\nabla u^{\varepsilon}$.

Then, passing to a (not relabeled) subsequence if needed, it follows that $u^{\varepsilon} \rightarrow u$ a.e. in $\Omega \times Q$.
Hence $u$ satisfies (6), or, in other words, $\partial_{t}\left[u-\int_{0}^{t} \Phi_{N} d W\right]-\Delta_{p(\cdot)} u=0$.
In particular, since $\Phi_{N}$ is regular, one gets that $u-\int_{0}^{t} \Phi_{N} d W \in \mathcal{E}$ with $\partial_{t}\left[u-\int_{0}^{t} \Phi_{N} d W\right] \in \mathcal{E}^{\prime}$.
We need now to prove that $u \in L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$. We already know that $u: \Omega \times Q \rightarrow L^{2}(D)$ is a stochastic process. Since $u(\omega, \cdot) \in W_{\omega}(Q) \hookrightarrow C\left([0, T] ; L^{2}(D)\right)$ for a.e. $\omega \in \Omega$, the measurability follows from [2, Prop. 3.17 p. 84$]$ with arguments as in [4, Cor. 1.1 .2 , p. 8]. Then, a.s. in $\Omega$, the equation satisfied by $u$ yields $\partial_{t} v-\Delta_{p(\cdot)} u=0$, so that, for almost every $t \in[0, T]$,

$$
\frac{1}{2} \frac{d}{d t}\|v(t)\|_{L^{2}(D)}^{2}+\int_{D}|\nabla u|^{p(\omega, t, x)-2} \nabla u \cdot \nabla v d x=0
$$

Since, $\omega$ a.s.,

$$
\sup _{t \in[0, T]}\|v(\omega, t, \cdot)\|_{L^{2}(D)}^{2} \leq\left\|u_{0}\right\|_{L^{2}(D)}^{2}+2 \int_{0}^{T} \int_{D} \frac{1}{p^{-}}|\nabla u|^{p(\omega, s, x)}+\frac{1}{\left(p^{\prime}\right)^{-}}\left|\int_{0}^{s} \nabla \Phi_{N} d W\right|^{p^{\prime}(\omega, s, x)} d x d s
$$

with a right-hand side in $L^{1}(\Omega)$, one gets that $u, v \in L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right.$.
Lemma 4.1. For any $m, n \in \mathbb{N}$ and nice processes $\Phi_{N, n}, \Phi_{N, m}$ let $u_{n}$ be the solution to ( $P, \Phi_{N, n}$ ) and $u_{m}$ be the solution to $\left(P, \Phi_{N, m}\right)$. There exist constants $K_{1}, K_{2} \geq 0$ depending on the Burkholder-Davies-Gundy inequality such that

$$
\begin{align*}
& E\left(\left\|u_{n}\right\|_{C\left([0, T] ; L^{2}(D)\right)}^{2}+\int_{Q}\left|\nabla u_{n}\right|^{p(\cdot)} d(t, x)\right) \leq K_{1} E\left(\int_{0}^{T}\left\|\Phi_{N, n}\right\|_{H S\left(L^{2}(D)\right)}^{2} d t+\left\|u_{0}\right\|_{L^{2}(D)}^{2}\right)  \tag{9}\\
& E\left(\left\|\left(u_{n}-u_{m}\right)\right\|_{C\left([0, T] ; L^{2}(D)\right)}^{2}\right)+E \int_{Q}\left(\left|\nabla u_{n}\right|^{p(\cdot)-2} \nabla u_{n}-\left|\nabla u_{m}\right|^{p(\cdot)-2} \nabla u_{m}\right) \cdot \nabla\left(u_{n}-u_{m}\right) d(t, x) \\
& \leq K_{2} E \int_{0}^{T}\left\|\Phi_{N, n}-\Phi_{N, m}\right\|_{H S\left(L^{2}(D)\right)}^{2} d t . \tag{10}
\end{align*}
$$

Proof. Using the Itô formula in (5), it follows that for all $t \in[0, T]$, a.s. in $\Omega$, we have

$$
\begin{align*}
& \left\|u_{n}^{\varepsilon}(t)\right\|_{L^{2}(D)}^{2}+2 \int_{0}^{t} \int_{D}\left|\nabla u_{n}^{\varepsilon}\right|^{p(\cdot)} d x d s \\
\leq & 2 \int_{0}^{t}\left(u_{n}^{\varepsilon}, \Phi_{N, n}(\cdot)\right)_{L^{2}(D)} d W+\int_{0}^{t}\left\|\Phi_{N, n}\right\|_{H S\left(L^{2}(D)\right)}^{2} d s+\left\|u_{0}\right\|_{L^{2}(D)}^{2} \tag{11}
\end{align*}
$$

or, by subtracting (5) with $\Phi_{N, m}$ from (5) with $\Phi_{N, n}$,

$$
\begin{align*}
& \left\|\left(u_{n}^{\varepsilon}-u_{m}^{\varepsilon}\right)(t)\right\|_{L^{2}(D)}^{2}+2 \int_{0}^{t} \int_{D}\left(\left|\nabla u_{n}^{\varepsilon}\right|^{p(\cdot)-2} \nabla u_{n}^{\varepsilon}-\left|\nabla u_{m}^{\varepsilon}\right|^{p(\cdot)-2} \nabla u_{m}^{\varepsilon}\right) \cdot \nabla\left(u_{n}^{\varepsilon}-u_{m}^{\varepsilon}\right) d x d s \\
\leq & 2 \int_{0}^{t}\left(u_{n}^{\varepsilon}-u_{m}^{\varepsilon},\left[\Phi_{N, n}-\Phi_{N, m}\right](\cdot)\right)_{L^{2}(D)} d W+\int_{0}^{t}\left\|\Phi_{N, n}-\Phi_{N, m}\right\|_{H S\left(L^{2}(D)\right)}^{2} d s . \tag{12}
\end{align*}
$$

Thus, by passing to the limit with $\varepsilon \rightarrow 0$, to the supremum over $t$ and then taking the expectation, it follows that ( $c \geq 0$ being a constant)

$$
\begin{align*}
& E\left(\sup _{t \in[0, T]}\left\|u_{n}(t)\right\|_{L^{2}(D)}^{2}\right)+E \int_{0}^{T} \int_{D}\left|\nabla u_{n}\right|^{p(\cdot)} d x d s \\
\leq & c E\left(\sup _{t \in[0, T]} \int_{0}^{t}\left(u_{n}, \Phi_{N, n}(\cdot)\right)_{L^{2}(D)} d W\right)+c\left\|\Phi_{N, n}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2}+c\left\|u_{0}\right\|_{L^{2}(D)}^{2}, \tag{13}
\end{align*}
$$

$$
\begin{align*}
& E\left(\sup _{t \in[0, T]}\left\|\left(u_{n}-u_{m}\right)(t)\right\|_{L^{2}(D)}^{2}\right)+E \int_{0}^{T} \int_{D}\left(\left|\nabla u_{n}\right|^{p(\cdot)-2} \nabla u_{n}-\left|\nabla u_{m}\right|^{p(\cdot)-2} \nabla u_{m}\right) \cdot \nabla\left(u_{n}-u_{m}\right) d x d s \\
\leq & c E\left(\sup _{t \in[0, T]} \int_{0}^{t}\left(u_{n}-u_{m},\left[\Phi_{N, n}-\Phi_{N, n}\right](\cdot)\right)_{L^{2}(D)} d W\right)+c\left\|\Phi_{N, n}-\Phi_{N, m}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2} \tag{14}
\end{align*}
$$

Using Burkholder, Hölder and Young inequalities on (13) we get for any $\gamma>0$

$$
\begin{equation*}
E\left(\sup _{t \in[0, T]} \int_{0}^{t}\left(u_{n}, \Phi_{N, n}(\cdot)\right)_{L^{2}(D)} d W\right) \leq 3 E\left(\int_{0}^{T}\left\|\left(u_{n}, \Phi_{N, n}(\cdot)\right)_{L^{2}(D)}\right\|_{H S\left(L^{2}(D), \mathbb{R}\right)}^{2} d s\right)^{1 / 2} \tag{15}
\end{equation*}
$$

where

$$
\begin{align*}
\left\|\left(u_{n}, \Phi_{N, n}(\cdot)\right)_{L^{2}(D)}\right\|_{H S\left(L^{2}(D), \mathbb{R}\right)}^{2} & =\sum_{k=1}^{\infty}\left|\left(u_{n}, \Phi_{N, n}\left(e_{k}\right)\right)_{L^{2}(D)}\right|^{2} \\
& =\sum_{k=1}^{\infty}\left|\left(\Phi_{N, n}^{*}\left(u_{n}\right), e_{k}\right)_{L^{2}(D)}\right|^{2} \leq\left\|\Phi_{N, n}\right\|_{H S\left(L^{2}(D)\right)}^{2}\left\|u_{n}\right\|_{L^{2}(D)}^{2} \tag{16}
\end{align*}
$$

and therefore

$$
\begin{aligned}
E\left(\sup _{t \in[0, T]} \int_{0}^{t}\left(u_{n}, \Phi_{N, n}\right)_{L^{2}(D)} d W\right) & \leq 3 E\left(\int_{0}^{T}\left\|\Phi_{N, n}\right\|_{H S\left(L^{2}(D)\right)}^{2}\left\|u_{n}\right\|_{L^{2}(D)}^{2} d t\right)^{1 / 2} \\
& \leq 3 E\left[\left(\sup _{t \in[0, T]}\left\|u_{n}\right\|_{L^{2}(D)}^{2}\right)^{1 / 2}\left(\int_{0}^{T}\left\|\Phi_{N, n}\right\|_{H S\left(L^{2}(D)\right)}^{2}\right)^{1 / 2}\right] \\
& \leq 3 \gamma E\left(\sup _{t \in[0, T]}\left\|u_{n}\right\|_{L^{2}(D)}^{2}\right)+\frac{3}{\gamma}\left\|\Phi_{N, n}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2}
\end{aligned}
$$

and similarly on (14),

$$
\begin{align*}
& E\left(\sup _{t \in[0, T]} \int_{0}^{t}\left(u_{n}-u_{m},\left[\Phi_{N, n}-\Phi_{N, m}\right](\cdot)\right)_{L^{2}(D)} d W\right) \\
\leq & 3 \gamma E\left(\sup _{t \in[0, T]}\left\|u_{n}-u_{m}\right\|_{L^{2}(D)}^{2}\right)+\frac{3}{\gamma}\left\|\Phi_{N, n}-\Phi_{N, m}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2} . \tag{17}
\end{align*}
$$

Plugging (15) into (13), (17) into (14) and choosing $\gamma>0$ small enough yield Lemma 4.1.
Remark 4.1. It is an open question if the Itô formula is directly available for a solution of (6) since we are not in Bochner spaces: The stochastic energy has to be defined in different Banach spaces depending on $t \in[0, T]$ and $\omega \in \Omega$. That is why we need to apply the Itô formula to $u^{\varepsilon}$, and then pass to the limit. But then, only an inequality is obtained.

### 4.2. Existence for arbitrary $\Phi_{N}$

Proposition 4.3. For any $N \in \mathbb{N}^{*}$ and any $\Phi_{N}$ defined as in the beginning of the section, there exists a unique solution $u$ to $\left(P, \Phi_{N}\right)$, i.e., $u \in \mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right) \cap N_{w}^{2}\left(0, T ; L^{2}(D)\right)\right.$ such that a.s.

$$
\begin{equation*}
u(t)-u_{0}-\int_{0}^{t} \Delta_{p(\cdot)} u d s=\int_{0}^{t} \Phi_{N} d W \tag{18}
\end{equation*}
$$

for all $t \in[0, T]$, a.e. in $D$.
Moreover, if $u_{1}, u_{2}$ are the solutions to $\left(P, \Phi_{N, 1}\right),\left(P, \Phi_{N, 2}\right)$ respectively, then:

$$
\begin{align*}
& E\left(\sup _{t}\left\|\left(u_{1}-u_{2}\right)(t)\right\|_{L^{2}(D)}^{2}\right)+E \int_{Q}\left(\left|\nabla u_{1}\right|^{p(\cdot)-2} \nabla u_{1}-\left|\nabla u_{2}\right|^{p(\cdot)-2} \nabla u_{2}\right) \cdot \nabla\left(u_{1}-u_{2}\right) d(t, x) \\
\leq & C E \int_{0}^{T}\left\|\Phi_{N, 1}-\Phi_{N, 2}\right\|_{H S\left(L^{2}(D)\right)}^{2} d t . \tag{19}
\end{align*}
$$

Proof. For any $k \in \mathbb{N}$ there exists a sequence $\left(\Phi_{n, k, N}\right)_{n} \subset S_{w}^{2}\left(0, T ; H_{0}^{j}(D)\right)$ converging to $\Phi_{N}\left(e_{k}\right)$ in $N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ when $n$ goes to $\infty$. If we define $\Phi_{n, N}: \Omega \times(0, T) \rightarrow H S\left(L^{2}(D)\right)$ by $\Phi_{n, N}\left(e_{k}\right)=\Phi_{n, k, N}$ for $k \leq N$, then $\Phi_{n, N}$ is a nice process such that $\Phi_{n, N} \rightarrow \Phi_{N}$ in $L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)$ when $n \rightarrow \infty$. Let $\left(u_{n}\right) \in \mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$ be the sequence of corresponding solutions to ( $P, \Phi_{n, N}$ ), then from (9) it follows that $\left(u_{n}\right)$ is a bounded sequence in $\mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right.$ ) and (10) ensures that ( $u_{n}$ ) is a Cauchy sequence in $L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right.$.

Hence there exists $u \in \mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$ such that $u_{n} \rightharpoonup u$ in $\mathcal{E}$ and $u_{n} \rightarrow u$ in $L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$.

Moreover there exists a full-measure set $\tilde{\Omega} \in \mathcal{F}$ such that, passing to a (not relabeled) subsequence if necessary, $u_{n} \rightarrow u$ in $C\left([0, T] ; L^{2}(D)\right)$ for all $\omega \in \tilde{\Omega}$. In particular, $u(0, \cdot)=u_{0}$ a.e. in $D$ for all $\omega \in \tilde{\Omega}$.

For $\mu=d(t, x) \otimes d P$ we have

$$
\int_{\Omega \times Q}\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)} d \mu=\int_{1<p<2}\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)} d \mu+\int_{p \geq 2}\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)} d \mu .
$$

Then, from (10) and the fundamental inequality ([6, Section 10$]$ ), for any $\xi, \eta \in \mathbb{R}^{d}$,

$$
\left(|\xi|^{p-2} \xi-|\eta|^{p-2} \eta\right) \cdot(\xi-\eta) \geq\left\{\begin{array}{l}
2^{2-p}|\xi-\eta|^{p}, p \geq 2 \\
(p-1)|\xi-\eta|^{2}\left(1+|\eta|^{2}+|\xi|^{2}\right)^{\frac{p-2}{2}}, 1 \leq p<2
\end{array}\right.
$$

it follows first that

$$
\begin{equation*}
\int_{p \geq 2}\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)} d \mu \leq 2^{p^{+}-2} K_{2}\left\|\Phi_{n, N}-\Phi_{m, N}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2} . \tag{20}
\end{equation*}
$$

Then, from the generalized Young inequality it follows secondly that, for any $0<\epsilon<1$,

$$
\int_{1<p<2}\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)} d \mu=\int_{1<p<2} \frac{\left|\nabla u_{n}-\nabla u_{m}\right|^{p(\cdot)}}{\left(1+\left|\nabla u_{n}\right|^{2}+\left|\nabla u_{m}\right|^{2}\right)^{p(\cdot) \frac{2-p(\cdot)}{4}}}\left(1+\left|\nabla u_{n}\right|^{2}+\left|\nabla u_{m}\right|^{2}\right)^{p(\cdot) \frac{2-p(\cdot)}{4}} d \mu
$$

$$
\begin{align*}
& \leq \int_{1<p<2} \epsilon^{\frac{p(\cdot)-2}{p(\cdot)}} \frac{\left|\nabla u_{n}-\nabla u_{m}\right|^{2}}{\left(1+\left|\nabla u_{n}\right|^{2}+\left|\nabla u_{m}\right|^{2}\right)^{\frac{2-p(\cdot)}{2}}} d \mu \\
& +\epsilon \int_{1<p<2}\left(1+\left|\nabla u_{n}\right|^{2}+\left|\nabla u_{m}\right|^{2}\right)^{\frac{p(\cdot)}{2}} d \mu \\
& \leq \frac{1}{\epsilon\left(p^{-}-1\right)} \int_{1<p<2}(p-1) \frac{\left|\nabla u_{n}-\nabla u_{m}\right|^{2}}{\left(1+\left|\nabla u_{n}\right|^{2}+\left|\nabla u_{m}\right|^{2}\right)^{\frac{2-p(\cdot)}{2}}} d \mu+K_{3} \epsilon \\
& \leq \frac{1}{\epsilon\left(p^{-}-1\right)} K_{2}\left\|\Phi_{n, N}-\Phi_{m, N}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2}+K_{3} \epsilon, \tag{21}
\end{align*}
$$

since the sequence $\left(u_{n}\right)$ is bounded in $L^{p(\cdot)}(\Omega \times Q)$ and $\mu$ is a finite measure.
From (20), (21) and $\lim _{n, m \rightarrow \infty}\left\|\Phi_{n, N}-\Phi_{m, N}\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2}=0$ one gets that $\nabla u_{n}$ is a Cauchy sequence in $L^{p(\cdot)}(\Omega \times Q)$, thus a converging sequence.

In conclusion, $u_{n}$ converges to $u$ in $\mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right) \cap N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ and, by a standard argument based on the Nemytskii operator induced by the Carathéodory function $G:(\omega, t, x, \xi) \in$ $\Omega \times Q \times \mathbb{R}^{d} \mapsto|\xi|^{p(\omega, t, x)-2} \xi \in \mathbb{R}^{d},\left|\nabla u_{n}\right|^{p(\cdot)-2} \nabla u_{n}$ converges to $|\nabla u|^{p(\cdot)-2} \nabla u$ in $L^{p^{\prime}(\cdot)}(\Omega \times Q)$ since $|G(\omega, t, x, \xi)|^{p^{\prime}(\omega, t, x)}=|\xi|^{p(\omega, t, x)}$.

Let us recall that, for any $n \in \mathbb{N}$, $u_{n}$ satisfies

$$
\begin{equation*}
\partial_{t}\left(u_{n}-\int_{0}^{t} \Phi_{n, N} d W\right)-\Delta_{p(\cdot)} u_{n}=0 \tag{22}
\end{equation*}
$$

in $\mathcal{E}^{\prime}$. Now we can choose a (not relabeled) subsequence of $\left(u_{n}\right)$ such that all previous convergence results hold true. For any test function $\phi(\omega, t, x)=\rho(\omega) \gamma(t) \nu(x)$ with $\rho \in L^{\infty}(\Omega), \gamma \in \mathcal{D}([0, T))$ and $\nu \in \mathcal{D}(D)$ we have

$$
\begin{align*}
& \left\langle\partial_{t}\left(u_{n}-\int_{0}^{t} \Phi_{n, N} d W\right), \phi\right\rangle_{\mathcal{E}^{\prime}, \mathcal{E}}=\int_{\Omega}\left\langle\partial_{t}\left(u_{n}-\int_{0}^{t} \Phi_{n, N} d W\right), \phi\right\rangle_{X_{\omega}^{\prime}, X_{\omega}} d P \\
= & -\int_{\Omega}\left\langle\left(u_{n}-\int_{0}^{t} \Phi_{n, N} d W\right), \partial_{t} \phi\right\rangle_{X_{\omega}^{\prime}, X_{\omega}} d P-\int_{\Omega \times D} u_{0} \varphi(\omega, 0, x) d x d P . \tag{23}
\end{align*}
$$

In particular $u_{n}$ satisfies

$$
\begin{equation*}
-\int_{\Omega \times Q}\left(u_{n}-\int_{0}^{t} \Phi_{n, N} d W\right) \cdot \partial_{t} \phi+\left|\nabla u_{n}\right|^{p(\cdot)-2} \nabla u_{n} \cdot \nabla \phi d \mu-\int_{\Omega \times D} u_{0} \varphi(\omega, 0, x) d x d P=0 \tag{24}
\end{equation*}
$$

for all $n \in \mathbb{N}$. When $n \rightarrow \infty$, from Itô isometry it follows that

$$
\begin{align*}
\int_{\Omega} \int_{0}^{T}\left\|\int_{0}^{t} \Phi_{n, N}-\Phi_{N} d W\right\|_{L^{2}(D)}^{2} d t d P & \leq \sup _{t \in[0, T]} E\left\|\int_{0}^{t} \Phi_{n, N}-\Phi_{N} d W\right\|_{L^{2}(D)}^{2} \\
& =E \int_{0}^{T}\left\|\Phi_{n, N}-\Phi_{N}\right\|_{H S\left(L^{2}(D)\right)}^{2} d t \rightarrow 0 \tag{25}
\end{align*}
$$

hence $\int_{0}^{\cdot} \Phi_{n, N} d W \rightarrow \int_{0}^{*} \Phi_{N} d W$ in $L^{2}\left(\Omega \times(0, T) ; L^{2}(D)\right)$ for $n \rightarrow \infty$. Therefore, using our convergence results, we are able to pass to the limit in (24) and obtain

$$
\begin{equation*}
\partial_{t}\left(u-\int_{0}^{t} \Phi_{N} d W\right)-\Delta_{p(\cdot)} u=0 \tag{26}
\end{equation*}
$$

in $\mathcal{E}^{\prime} .(26)$, and a classical argument of separability, imply that a.s.

$$
\begin{equation*}
\partial_{t}\left(u-\int_{0}^{t} \Phi_{N} d W\right)=\Delta_{p(\cdot)} u, \text { in } X_{\omega}^{\prime}(Q) \hookrightarrow L^{\alpha^{\prime}}\left(0, T ; W^{-1, \alpha^{\prime}}(D)\right) \tag{27}
\end{equation*}
$$

with $\alpha \geq p^{+}+2$. Moreover, a.s.

$$
u-\int_{0}^{t} \Phi_{N} d W \in C\left([0, T] ; L^{2}(D)\right)
$$

Thus we can integrate (27) to obtain a.s.

$$
\begin{equation*}
u(t)-u_{0}-\int_{0}^{t} \Delta_{p(\cdot)} u d s=\int_{0}^{t} \Phi_{N} d W \tag{28}
\end{equation*}
$$

in $L^{2}(D)$ for all $t \in[0, T]$.
To finish the proof, note that passing to the limit in (10) yields the stability inequality (19), and if we assume that $u_{1}, u_{2} \in \mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right) \cap N_{W}^{2}\left(0, T ; L^{2}(D)\right)$ are both satisfying (18), it follows that a.s. in $\Omega$

$$
\begin{equation*}
\partial_{t}\left(u_{1}-u_{2}\right)-\left(\Delta_{p(\cdot)} u_{1}-\Delta_{p(\cdot)} u_{2}\right)=0 \text { in }\left(X_{\omega}(Q)\right)^{\prime} \tag{29}
\end{equation*}
$$

Using $u_{1}-u_{2}$ as a test function in (29) and the integration by parts formula in $W_{\omega}(Q)$, we obtain the result of uniqueness.

### 4.3. Existence for cylindrical Wiener process

Now, our aim is to pass to the limit when $N \rightarrow \infty$ in Proposition 4.3. For $M, N \in \mathbb{N}^{*}$, let $u_{N}$ and $u_{M}$ be the solutions obtained by Proposition 4.3 corresponding to right-hand sides $\Phi_{N}$ and $\Phi_{M}$ respectively. Thanks to the stability result of the proposition, one has that

$$
\begin{align*}
E\left(\left\|\left(u_{N}-u_{M}\right)\right\|_{C\left([0, T] ; L^{2}(D)\right)}^{2}\right) & +E \int_{Q}\left(\left|\nabla u_{N}\right|^{p(\cdot)-2} \nabla u_{N}-\left|\nabla u_{M}\right|^{p(\cdot)-2} \nabla u_{M}\right) \cdot \nabla\left(u_{N}-u_{M}\right) d(t, x) \\
& \leq K_{2} E \int_{0}^{T}\left\|\Phi_{N}-\Phi_{M}\right\|_{H S\left(L^{2}(D)\right)}^{2} d t \tag{30}
\end{align*}
$$

Since $\Phi_{N}$ converges to $\Phi$ in $L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right.$ ), the right-hand side of (30) converges to 0 when $M, N \rightarrow \infty$. Therefore $\left(u_{N}\right)$ is a Cauchy sequence in $L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right)$. Using (20) and (21), we find that $\left(\nabla u_{N}\right)$ is a Cauchy sequence in $L^{p(\cdot)}(\Omega \times Q)^{d}$. Thus, there exists $u$ in $\mathcal{E} \cap L^{2}\left(\Omega ; C\left([0, T] ; L^{2}(D)\right)\right) \cap$
$N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ such that $u_{N} \rightarrow u$ in $\mathcal{E}$. Moreover, $\left|\nabla u_{N}\right|^{p(\cdot)-2} \nabla u_{N}$ converges to $|\nabla u|^{p(\cdot)-2} \nabla u$ in $L^{p^{\prime}(\cdot)}(\Omega \times Q)$ by a standard argument based again on the Nemytskii operator induced by the Carathéodory function $G:(\omega, t, x, \xi) \in \Omega \times Q \times \mathbb{R}^{d} \mapsto|\xi|^{p(\omega, t, x)-2} \xi \in \mathbb{R}^{d}$, since $|G(\omega, t, x, \xi)|^{p^{\prime}(\omega, t, x)}=|\xi|^{p(\omega, t, x)}$. In addition, from Itô isometry and the convergence of $\Phi_{N}$ to $\Phi$ in $L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)$ when $N \rightarrow \infty$ it follows that

$$
\int_{0}^{\dot{j}} \Phi_{N} d W \rightarrow \int_{0} \Phi d W
$$

in $L^{2}\left(\Omega \times(0, T) ; L^{2}(D)\right)$ when $N \rightarrow \infty$. From

$$
\begin{equation*}
-\int_{\Omega \times Q}\left(u_{N}-\int_{0}^{t} \Phi_{N} d W\right) \cdot \partial_{t} \phi+\left|\nabla u_{N}\right|^{p(\cdot)-2} \nabla u_{N} \cdot \nabla \phi d \mu-\int_{\Omega \times D} u_{0} \varphi(\omega, 0, x) d x d P=0 \tag{31}
\end{equation*}
$$

and with analogous arguments as in (26) and (27) it follows that

$$
\partial_{t}\left(u-\int_{0}^{t} \Phi d W\right) \in L^{\alpha^{\prime}}\left(0, T ; W^{-1, p^{\prime}}(D)\right), \quad u-\int_{0}^{t} \Phi d W \in C\left([0, T] ; L^{2}(D)\right)
$$

a.s. in $\Omega$ for $\alpha \geq 2+p^{+}$. Thus we can integrate (27) to obtain a.s.

$$
\begin{equation*}
u(t)-u_{0}-\int_{0}^{t} \Delta_{p(\cdot)} u d s=\int_{0}^{t} \Phi d W \tag{32}
\end{equation*}
$$

in $L^{2}(D)$ for all $t \in[0, T]$.
As a conclusion, on has that

Proposition 4.4. Proposition 4.3 still holds for any $\Phi \in N_{w}^{2}\left(0, T ; H S\left(L^{2}(D)\right)\right)$.

### 4.4. The multiplicative case

We want to apply Banach's fixed point theorem to the map

$$
\Psi: S \in N_{w}^{2}\left(0, T ; L^{2}(D)\right) \rightarrow u_{S} \in N_{w}^{2}\left(0, T ; L^{2}(D)\right)
$$

(see also [1]) where $u_{S}$ is the solution to $(P, H(\cdot, S))$ to deduce the existence of a unique solution $u$ of $(P, H)$ in the sense of Definition 3.1.

Let us first note that thanks to the assumptions on $\left(h_{k}\right)$ and by classical arguments based on Nemytskii operators, one has that $H(\cdot, S) \in N_{w}^{2}\left(0, T ; H S\left(L^{2}(D)\right)\right)$ when $S \in N_{w}^{2}\left(0, T ; L^{2}(D)\right)$.

Thus, thanks to Proposition 4.4, the mapping $\Psi$ is well-defined and $u$ is a solution to the multiplicative problem iff it is a fixed point for $\Psi$.

Set $S_{1}, S_{2} \in N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ and denote by $u_{1}=\Psi\left(S_{1}\right)$ and $u_{2}=\Psi\left(S_{2}\right)$. Thanks to the stability inequality for the additive case and to $(H 1)$ it follows that

$$
E\left\|\left(u_{1}-u_{2}\right)\right\|_{C\left([0, T] ; L^{2}(D)\right)}^{2} \leq C\left\|H\left(\cdot, S_{1}\right)-H\left(\cdot, S_{2}\right)\right\|_{L^{2}\left(\Omega \times(0, T) ; H S\left(L^{2}(D)\right)\right)}^{2}
$$

$$
\begin{align*}
& =C E \int_{0}^{T} \int_{D} \sum_{k=1}^{\infty}\left|h_{k}\left(\omega, t, S_{1}(\omega, t, x)\right)-h_{k}\left(\omega, t, S_{2}(\omega, t, x)\right)\right|^{2} d x d t \\
& \leq C\left\|S_{1}-S_{2}\right\|_{N_{w}^{2}\left(0, T ; L^{2}(D)\right)}^{2} \tag{33}
\end{align*}
$$

where $C \geq 0$ is a constant that may change from line to line. From (33) it follows that

$$
\begin{equation*}
\int_{0}^{T} E\left\|\left(u_{1}-u_{1}\right)(t)\right\|_{L^{2}(D)}^{2} e^{-\alpha t} d t \leq \frac{C}{\alpha}\left(1-e^{-\alpha T}\right) \int_{0}^{T} E\left\|S_{1}-S_{2}\right\|_{L^{2}(D)}^{2} e^{-\alpha t} d t \tag{34}
\end{equation*}
$$

for $\alpha>0$ and therefore the mapping $S$ has a fixed point in $N_{w}^{2}\left(0, T ; L^{2}(D)\right)$ for a suitable value of $\alpha$.
This finishes the proof of Theorem 3.1.

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