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## Stochastic differential equations in Banach spaces, variational formulation

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Analisi matematica. - Stochastic differential equations in Banach spaces, variational formulation ${ }^{(*)}$. Nota ${ }^{(* *)}$ di Giuseppe Da Prato, Mimmo Iannelli e Luciano Tubaro, presentata dal Corrisp. G. Stampacchia.

Riassunto. - Si danno risultati di esistenza e unicità, da un punto di vista variazionale, della soluzione per una equazione differenziale stocastica in spazi di Hilbert, in condizioni di non-lipschitzianità.

## I. Introduction

Let us recall the definition of abstract Wiener space: let $H_{0}$ be a separable Hilbert space; we will denote the set of finite projections on $\mathrm{H}_{0}$ by $\mathscr{F}$. Let X be a Banach space of which $\mathrm{H}_{0}$ is a dense subspace.

Definition i. ( $\mathrm{H}_{\mathbf{0}}, \mathrm{X}$ ) is an abstract Wiener space if for any $\varepsilon>0$ there exists a projection $\mathrm{P}_{\varepsilon} \in \mathscr{F}$ such that

$$
P \perp P_{\varepsilon} \Rightarrow \mu\left(\left\{x:|P x|_{X}>\varepsilon\right\}\right)<\varepsilon
$$

where $\mathrm{P} \in \mathscr{F}$ and

$$
\mu\left(\left\{x:|\mathrm{P} x|_{\mathrm{X}}>\varepsilon\right\}\right)=(2 \pi)^{-\frac{n}{2}} \int_{\left\{x \in \mathrm{P}\left(\mathrm{H}_{0}\right):|\mathrm{P} x|_{\mathrm{X}}>\varepsilon\right\}} e^{-\frac{|x|_{\mathrm{H}_{0}}^{2}}{2}} \mathrm{~d} x
$$

where $n=\operatorname{dim} \mathrm{P}\left(\mathrm{H}_{0}\right)$.
Theorem I. Let $\left(\mathrm{H}_{\mathbf{0}}, \mathrm{X}\right)$ an abstract Wiener space; if B is any Borel set in $\mathrm{P}\left(\mathrm{H}_{0}\right)$, where $\mathrm{P} \in \mathscr{F}$, the measure

$$
\mu_{t}\left(\mathrm{P}^{-1}(\mathrm{~B})\right)=(2 \pi t)^{-\frac{n}{2}} \int_{\mathrm{B} \cap \mathrm{P}\left(\mathrm{H}_{0}\right)} e^{-\frac{|x|_{\mathrm{H}_{0}}^{2}}{2 t}} \mathrm{~d} x \quad n=\operatorname{dim} \mathrm{P}\left(\mathrm{H}_{1}\right)
$$

defined on all cylindrical sets ${ }^{(1)}$ in $\mathrm{H}_{0}$, as an extension to cylindrical sets of X in such a way it is countably additive.

Example. Let X be a Hilbert space, $\mathrm{H}_{0}=\sqrt{\mathrm{S}} \mathrm{X}$ where S is a strictly positive trace-operator on X .
(*) Lavoro eseguito nell'ambito del G.N.A.F.A del C.N.R.
(**) Pervenuta all'Accademia il 7 settembre 1976.
(1) A cylindrical set is any set like $\mathrm{P}^{-1}(\mathrm{~B})$ for some $\mathrm{P} \in \mathscr{F}$ and some B , Borel set in $P\left(H_{0}\right)$.

Let V be a separable reflexive Banach space and H a Hilbert space such that $\mathrm{V} \leftrightarrows \mathrm{H} \subset \mathrm{V}^{\prime}{ }^{(2)}$. Let ( $\Omega, \mathscr{E}, \mathrm{P}$ ) be a probability space, $\mathrm{H}_{0} \subset \mathrm{~V}$ a Hilbert space such that $\left(\mathrm{H}_{0}, \mathrm{~V}\right)$ is an abstract Wiener space, $\mathrm{W}(t)$ a Wiener process in $\left(\mathrm{H}_{0}, \mathrm{~V}\right)$ and $\mathscr{F}_{t}$ the smallest $\sigma$-algebra such that $\mathrm{W}(s)$, $0 \leq s \leq t$, is measurable ${ }^{(3)}$.

Given the mappings

$$
\begin{aligned}
& f: \mathrm{V} \rightarrow \mathrm{~V}^{\prime} \\
& \mathrm{G}: \mathrm{V} \rightarrow \mathrm{~L}_{\lambda}^{*}\left(\mathrm{H}_{0}\right)
\end{aligned}
$$

where $L_{\lambda}^{*}\left(H_{0}\right)$ is the space of all operators $T$ of the form $T=\lambda I+T_{0}$, with $\lambda$ fixed in $\mathbf{R}$ and $T_{0}$ variable in the space of Hilbert-Schmidt operators in $H_{0}$. Putting $L^{*}\left(H_{0}\right)=\bigcup_{\lambda \in \mathbf{R}} L_{\lambda}^{*}\left(H_{0}\right)$, it is easy to see that $L^{*}\left(H_{0}\right)$ is a Hilbert space ${ }^{(4)}$; our G's are particular functions from V in $L^{*}\left(H_{0}\right)$. We will study the stochastic Cauchy problem:

$$
\left\{\begin{array}{l}
\mathrm{d} u(t)=f(u(t)) \mathrm{d} t+\mathrm{G}(u(t)) \mathrm{d} \mathrm{~W}_{t}  \tag{P}\\
u(0)=u_{0} \in \mathrm{~L}_{0}^{p}(\mathrm{~V})
\end{array}\right.
$$

with $p \geq 2$. In the non-stochastic case ( $G=0$ ), if $-f$ is assumed to be a monotone, hemicontinuous and coercive mapping, a solution of ( P ) is found by the Faedo-Galerkin approximation and the same method works in the case $G=I$ (see [I]). In the general case it is not possible to use this procedure for there is no general existence theorem for finite-dimensional stochastic equations and hence it is not possible to get approximate solutions of ( P ) without adding some other condition ${ }^{(6)}$.

In this paper we show existence for ( P ) via the Yosida approximation in H putting on $f$ and G suitable conditions that in the non-stochastic case reduce to the classical hypothesis used in a variational framework. These conditions have been also used in [8], together with some additional assumptions which allow to use the Faedo-Galerkin approximation.
(2) $\mathrm{X} \subset \rightarrow \mathrm{Y}$, where $\mathrm{X}, \mathrm{Y}$ are two B -spaces, means: $\mathrm{X} \subset \mathrm{Y}$ with continuous injections.
(3) More generally $\mathscr{F}_{t}$ is an one-parameter family of $\sigma$-algebras such that $\mathrm{W}_{s}$ is $\mathscr{F}_{t}$-measurable for $0 \leq s \leq t$ and, for $t \geq s, \mathscr{F}_{s}$ and $\mathrm{W}_{t}-\mathrm{W}_{s}$ are independent (see [4]).
(4) $L^{*}\left(\mathrm{H}_{0}\right)$ is a Hilbert space with respect to scalar product

$$
(\mathrm{T}, \mathrm{~S})_{*}=\sum_{v=1}^{\infty}\left(\mathrm{T} f_{v}, \mathrm{~S} f_{v}\right)
$$

where $\left\{f_{v}\right\}$ is any orthonormal basis in $\mathrm{H}_{0}$; the scalar product is independent from the choice of the basis.
(5) $\mathrm{L}_{\boldsymbol{t}}^{\boldsymbol{p}}(\mathrm{X})$ is the space of functions $\mathrm{L}^{p}\left(\Omega, \mathscr{F}_{t}, \mathrm{X}\right)$, where X is a Banach space.
(6) Concerning existence in the finite dimensional case, see [3].

We look for a solution of $(\mathrm{P})$ with the following properties:

$$
\begin{equation*}
u \in \mathrm{C}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) \tag{I}
\end{equation*}
$$

$$
\begin{equation*}
f(u) \in \mathrm{L}^{q}\left(0, \mathrm{~T} ; \mathrm{L}_{l}^{q}\left(\mathrm{~V}^{\prime}\right)\right) \quad \frac{\mathrm{I}}{p}+\frac{\mathrm{I}}{q}=\mathrm{I} \tag{2}
\end{equation*}
$$

$$
\begin{equation*}
\mathrm{G}(u) \in \mathrm{L}^{2}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{2}\left(\mathrm{~L}^{*}\left(\mathrm{H}_{0}\right)\right)\right) \tag{3}
\end{equation*}
$$

so that the stochastic differential in ( P ) makes sense.
The following lemmas will be used to get estimates, actually they precise our use of the Itô formula.

Lemma 1 . Let $u \in \mathrm{C}\left(\mathrm{o}, \mathrm{T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right)$ have the stochastic differential

$$
\mathrm{d} u=a(t) \mathrm{d} t+\mathrm{B}(t) \mathrm{dW}_{t}
$$

and suppose that $a \in \mathrm{~L}^{2}\left(\mathrm{o}, \mathrm{T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right), \mathrm{B} \in \mathrm{L}^{2}\left(\mathrm{o}, \mathrm{T} ; \mathrm{L}_{t}^{2}\left(\mathrm{~L}^{*}\left(\mathrm{H}_{0}\right)\right)\right.$ ); then:
(4) $\quad|u(t)|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}+\int_{0}^{t} \mathrm{E}\left\{2(a(s), u(s))+\|\mathrm{B}(s)\|_{\mathrm{L}^{*}\left(\mathrm{H}_{0}\right)}^{2}\right\} \mathrm{d} s$.

Proof. Let $\varphi_{k}: \mathbf{R} \rightarrow \mathbf{R}$ be a sequence such that

$$
\left\{\begin{array}{lll}
\varphi_{k} \in \mathrm{C}^{2}(\mathbf{R}) \quad, \quad \varphi_{k}(0)=0 \quad, \quad\left|\varphi_{k}(r)\right| \leq c k \quad, \quad \varphi_{k}(r) \rightarrow r \\
\left|\varphi_{k}^{\prime}(r)\right| \leq c & , \quad \varphi_{k}^{\prime}(r) \rightarrow \mathrm{I} \quad, \quad\left|\varphi_{k}^{\prime \prime}(r)\right| \leq \frac{c}{k}
\end{array}\right.
$$

Define the mapping $\Phi_{k}: H \rightarrow V$ putting:

$$
\Phi_{k}(x)=\sum_{i}^{k} \varphi_{k}\left(x_{i}\right) e_{i}
$$

$$
\forall x \in \mathrm{H}
$$

where $\left\{e_{i}\right\} \subset \mathrm{V}$ is a basis in H and $x_{i}=\left(x, e_{i}\right)$. Then we have:

$$
\begin{aligned}
\mathrm{d}\left|\Phi_{k}(u(t))\right|^{2} & =\left[2\left(\Phi_{k}(u(t)), \Phi_{k}^{\prime}[u(t)] a(t)\right)+\right. \\
& +\left\|\Phi_{k}^{\prime}[u(t)] \mathrm{B}(t)\right\|_{\mathrm{L}^{*}\left(\mathrm{H}_{0}\right)}^{2}+ \\
& \left.+\left(\Phi_{k}(u(t)), \mathrm{TR}\left\{\Phi_{k}^{\prime \prime}[u(t)](\mathrm{B}(t) \cdot, \mathrm{B}(t) \cdot)\right\}\right)\right] \mathrm{d} t+ \\
& +2\left(\Phi_{k}(u(t)), \Phi_{k}^{\prime}[u(t)] \mathrm{B}(t) \mathrm{dW}_{t}\right) .
\end{aligned}
$$

Hence:

$$
\begin{aligned}
& \left|\Phi_{k}(u(t))\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|\Phi_{k}\left(u_{0}\right)\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}+ \\
& +\int_{0}^{t} \mathrm{E}\left\{2\left(\Phi_{k}(u(s)), \Phi_{k}^{\prime}[u(s)] a(s)\right)+\left\|\Phi_{k}^{\prime}[u(s)] \mathrm{B}(s)\right\|_{\mathrm{L}^{*}\left(\mathrm{H}_{0}\right)}^{2}+\right. \\
& \left.+\left(\Phi_{k}(u(s)), \operatorname{TR}\left\{\Phi_{k}^{\prime \prime}[u(s)](\mathrm{B}(s) \cdot, \mathrm{B}(s) \cdot)\right\}\right)\right\} \mathrm{d} s
\end{aligned}
$$

and going to the limit we get (4).

Lemma 2. Let us suppose that

$$
\begin{aligned}
& u \in \mathrm{~L}^{p}\left(0, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) \quad, \quad u(0)=u_{0} \in \mathrm{~L}_{0}^{p}(\mathrm{~V}) \\
& a \in \mathrm{~L}^{q}\left(0, \mathrm{~T} ; \mathrm{L}_{t}^{q}\left(\mathrm{~V}^{\prime}\right)\right) \\
& \mathrm{A} \in \mathrm{~L}^{p}\left(0, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) \quad \text { where } \mathrm{A}(t)=\int_{0}^{t} a(s) \mathrm{d} s \\
& \mathrm{~B} \in \mathrm{~L}^{2}\left(0, \mathrm{~T} ; \mathrm{L}_{t}^{2}\left(\mathrm{~L}^{*}\left(\mathrm{H}_{0}\right)\right)\right)
\end{aligned}
$$

such that
then

$$
u(t)=u_{0}+\int_{0}^{t} a(s) \mathrm{d} s+\int_{0}^{t} \mathrm{~B}(s) \mathrm{dW}_{s}
$$

(4) $\quad|u(t)|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}+\int_{0}^{t} \mathrm{E}\left\{2|a(s), u(s)|+\|\mathrm{B}(s)\|_{L^{*}\left(\mathrm{H}_{0}\right)}^{2}\right\} \mathrm{d} s$.

Proof. We state that it is possible to find a sequence $\mathrm{A}_{n} \in \mathscr{D}\left([\mathrm{o}, \mathrm{T}] ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right)$ such that

$$
\begin{array}{ll}
\mathrm{A}_{n} \rightarrow \mathrm{~A} & \text { in } \mathrm{L}^{p}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) \\
\mathrm{A}_{n}^{\prime} \rightarrow a & \text { in } \mathrm{L}^{q}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{q}\left(\mathrm{~V}^{\prime}\right)\right) .
\end{array}
$$

Then consider

$$
u_{n}(t)=u_{0}+\mathrm{A}_{n}(t)+\int_{0}^{t} \mathrm{~B}(s) \mathrm{dW}_{s} .
$$

Clearly $u_{n} \rightarrow u$ in $L^{p}\left(o, T ; L_{t}^{p}(\mathrm{~V})\right)$, hence in $\mathrm{L}^{2}\left(\mathrm{o}, \mathrm{T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right)$; besides we can apply Lemma 1 to $\left(4^{\prime \prime}\right)$ to get

$$
\left|u_{n}(t)\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}+2 \int_{0}^{t} \mathrm{E}\left\{\left(\mathrm{~A}_{n}^{\prime}(s), u_{n}(s)\right)\right\} \mathrm{d} s+\int_{0}^{t} \mathrm{E}\left\{\|\mathrm{~B}(s)\|_{\mathrm{L}^{*}\left(\mathrm{H}_{0}\right)}^{2}\right\} \mathrm{d} s
$$

from which we get ( $4^{\prime}$ ) as $n \rightarrow \infty$. At last our statement at the beginning of the proof can be proved by adapting the proof of Theorem 2.I, Chapter I, of [7].

## 2. Dissipativity, Coercivity and the approximate problem

Let us consider the following assumption:

$$
\begin{equation*}
2\langle f(x)-f(y), x-y\rangle+\|\mathrm{G}(x)-\mathrm{G}(y)\|_{\mathrm{L}}^{2}{ }^{*}\left(\mathrm{H}_{0}\right) \leq 0 \quad \forall x, y \in \mathrm{~V} . \tag{1}
\end{equation*}
$$

The pair ( $f, G$ ) will be said decreasing if $\left(\mathrm{H}_{1}\right)$ is verified. We remark
that if $\left(\mathrm{H}_{1}\right)$ is verified then the mapping $-f: \mathrm{V} \rightarrow \mathrm{V}^{\prime}$ is monotone. The first consequence of $\left(\mathrm{H}_{1}\right)$ is the following theorem:

ThEOREM 2. Assume that $\left(\mathrm{H}_{1}\right)$ is verified, and let $u$ and $v$ be solutions of $(\mathrm{P})$ with initial data $u_{0}$ and $v_{0}$ respectively; then the following estimate is true:

$$
|u(t)-v(t)|_{L_{l}^{2}(\mathrm{H})}^{2} \leq\left|u_{0}-v_{0}\right|_{L_{0}(\mathrm{H})}^{2} .
$$

The proof of Theorem 2 is got by Lemma 2 and it means uniqueness of the solution of ( P ).

Our second assumption is

$$
\begin{equation*}
2\langle f(x), x\rangle+\|\mathrm{G}(x)\|_{\mathrm{L}^{*}\left(\mathrm{H}_{0}\right)}^{2} \leq-\omega\|x\|^{p} \quad \forall x \in \mathrm{~V} \tag{2}
\end{equation*}
$$

where $\omega>0$. If $\left(\mathrm{H}_{2}\right)$ is verified then the pair $(f, \mathrm{G})$ will be said to be coercive; obviously $\left(\mathrm{H}_{2}\right)$ yields coercivity for the mapping - $f$.

To complete the picture we also consider the following assumptions on $f$ :
$\left(\mathrm{H}_{3}\right) \quad f$ is hemicontinuous and $\quad\|f(x)\|_{\mathrm{v}^{\prime}} \leq k\|x\|^{p-1}$.
In the following we suppose that $\left(\mathrm{H}_{1}\right),\left(\mathrm{H}_{2}\right),\left(\mathrm{H}_{3}\right)$ are verified to define the approximate problems and show convergence to a solution of problem ( P ).

First of all we consider the following mapping:

$$
\tilde{f}:\left\{\begin{align*}
\mathrm{D}_{\tilde{f}} & =\{x \in \mathrm{~V} \mid f(x) \in \mathrm{H}\}  \tag{5}\\
\tilde{f}(x) & =f(x) \quad \forall x \in \mathrm{D}_{\tilde{f}}
\end{align*}\right.
$$

Owing to the assumptions $\tilde{f}$ is a maximal dissipative operator in H so that we can define the Yosida operators

$$
\begin{equation*}
\mathrm{J}_{n}=\left(\mathrm{I}-\frac{\mathrm{I}}{n} \tilde{f}\right)^{-1}: \mathrm{H} \rightarrow \mathrm{~V} \quad n>0 \tag{6}
\end{equation*}
$$

$$
\begin{equation*}
f_{n}=f \circ \mathrm{~J}_{n}=n\left(\mathrm{~J}_{n}-\mathrm{I}\right): \quad \mathrm{H} \rightarrow \mathrm{H} \quad n>0 \tag{7}
\end{equation*}
$$

with the well known properties:

$$
\begin{equation*}
\left|\mathrm{J}_{n}\right|_{\mathrm{L}} \leq \mathrm{I} \quad ; \quad\left|f_{n}\right|_{\mathrm{L}} \leq 2 n \tag{8}
\end{equation*}
$$

Yet we define:

$$
\begin{equation*}
\mathrm{G}_{n}=\mathrm{G} \circ \mathrm{~J}_{n}: \mathrm{H} \rightarrow \mathrm{~L}^{*}\left(\mathrm{H}_{0}\right) \quad n>0 . \tag{9}
\end{equation*}
$$

From $\left(\mathrm{H}_{1}\right)$ is follows that $\mathrm{G}_{n}$ is Lipschitz continuous so that the approximate problem:

$$
\left\{\begin{array}{l}
\mathrm{d} u_{n}(t)=f_{n}\left(u_{n}(t)\right) \mathrm{d} t+\mathrm{G}_{n}\left(u_{n}(t)\right) \mathrm{d} W_{t}  \tag{n}\\
u_{n}(\mathrm{o})=u_{0}
\end{array}\right.
$$

has a unique solution $u_{n} \in \mathrm{C}\left(0, \mathrm{~T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right)$. In the next section we state some estimates on the sequence of approximate solutions, hence existence of one solution of ( P ).

## 3. Existence for (P)

Let $\left\{u_{n}\right\}$ be the sequence of solutions of $\left(\mathrm{P}_{n}\right)$. We have first:
Proposition I." Let $\left(\mathrm{H}_{1}\right),\left(\mathrm{H}_{2}\right),\left(\mathrm{H}_{3}\right)$ be verified; then:

Proof. It is only worth proving (io) and (11), as (12) and (13) easily follow from these. Now from ( $\mathrm{P}_{n}$ ) it follows (see Lemma 1 ):

$$
\begin{array}{ll}
\text { (10) } & u_{n} \text { is a bounded sequence in } \mathrm{C}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right) \\
\text { (II) } & \mathrm{J}_{n} u_{n} \text { is a bounded sequence in } \mathrm{L}^{p}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) \\
\text { (I2) } & f_{n} u_{n} \text { is a bounded sequence in } \mathrm{L}^{q}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{q}\left(\mathrm{~V}^{\prime}\right)\right) \\
\text { (I3) } & \mathrm{G}_{n} u_{n} \text { is a bounded sequence in } \mathrm{L}^{2}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{2}\left(\mathrm{~L}^{*}\left(\mathrm{H}_{0}\right)\right)\right) . \tag{II}
\end{array}
$$

$$
\begin{equation*}
\left|u_{n}(t)\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2^{2}}+ \tag{14}
\end{equation*}
$$

$$
+\int_{0}^{t} \mathrm{E}\left\{2\left\langle f_{n}\left(u_{n}(s)\right), u_{n}(s)\right\rangle\left\|\mathrm{G}_{n}\left(u_{n}(s)\right)\right\|_{L^{*}\left(\mathrm{H}_{0}\right)}^{2}\right\} \mathrm{d} s
$$

Let us remember that

$$
\left\langle f_{n}\left(u_{n}(s)\right), u_{n}(s)\right\rangle=\left\langle f_{n}\left(u_{n}(s)\right), \mathrm{J}_{n}\left(u_{n}(s)\right)\right\rangle-\frac{1}{n}\left|f_{n}\left(u_{n}(s)\right)\right|^{2}
$$

from which, because of the coercivity,

$$
\left|u_{n}(t)\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2} \leq-\omega \int_{0}^{t}\left\|\mathrm{~J}_{n}\left(u_{n}(s)\right)\right\|_{\mathrm{L}^{p}(\mathrm{~V})}^{p} \mathrm{~d} s+\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}
$$

We finally have:
Theorem 3. Let $\left(\mathrm{H}_{1}\right),\left(\mathrm{H}_{2}\right),\left(\mathrm{H}_{3}\right)$ be verified; then there exists at least one solution of problem (P).

Proof. Let us pick from the sequence $\left\{u_{n}\right\}$ a subsequence ${ }^{(7)}$ such that

$$
\begin{array}{rll}
u_{n} \rightarrow u & \text { in } & \mathrm{L}^{\infty}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{2}(\mathrm{H})\right) \\
\mathrm{J}_{n} u_{n} \rightarrow v & \text { in } \mathrm{L}^{p}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{p}(\mathrm{~V})\right) & \text { weak } \\
f_{n} u_{n} \rightarrow \chi & \text { in } & \mathrm{L}^{q}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{q}\left(\mathrm{~V}^{\prime}\right)\right) \\
\mathrm{G}_{n} u_{n} \rightarrow \psi & \text { in } \mathrm{L}^{2}\left(\mathrm{o}, \mathrm{~T} ; \mathrm{L}_{t}^{2}\left(\mathrm{~L}^{*}\left(\mathrm{H}_{0}\right)\right)\right) & \text { weak }  \tag{18}\\
\text { weak } .
\end{array}
$$

(7) We will denote such a subsequence $\left\{u_{n}\right\}$ again.

Clearly (17), (16) and (15) imply that $u=v$ and that:

$$
\begin{equation*}
u(t)=u_{0}+\int_{0}^{t} \chi(s) \mathrm{d} s+\int_{0}^{t} \psi(s) \mathrm{dW}_{s} \tag{19}
\end{equation*}
$$

We have to show that:

$$
\chi(s)=f(u(s)) \quad, \quad \psi(s)=\mathrm{G}(u(s))
$$

which will be done adapting a classical method in abstract evolution equations (see [6]).

Let us consider

$$
\begin{aligned}
\mathrm{X}_{n}=\int_{0}^{\mathrm{T}} \mathrm{E} & \left\{2\left(f_{n}\left(u_{n}(s)\right)-f(v(s)), \mathrm{J}_{n}\left(u_{n}(s)\right)-v(s)\right)\right\} \mathrm{d} s+ \\
& +\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left\|\mathrm{G}_{n}\left(u_{n}(s)\right)-\mathrm{G}(v(s))\right\|_{L^{*}\left(\mathrm{H}_{0}\right)}^{2}\right\} \mathrm{d} s
\end{aligned}
$$

It is $\mathrm{X}_{n} \leq 0$ because of the dissipativity, on the other side:

$$
\begin{aligned}
\mathrm{X}_{n} & =2 \int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left(f_{n}\left(u_{n}(s)\right), \mathrm{J}_{n}\left(u_{n}(s)\right)\right)\right\} \mathrm{d} s+\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left\|\mathrm{G}_{n}\left(u_{n}(s)\right)\right\|_{*}^{2}\right\} \mathrm{d} s+ \\
& -2 \int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left(f(v(s)), \mathrm{J}_{n}\left(u_{n}(s)\right)-v(s)\right)\right\} \mathrm{d} s+ \\
& -\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left(\mathrm{G}(v(s)), \mathrm{G}_{n}\left(u_{n}(s)\right)-\mathrm{G}(v(s))\right)_{*}\right\} \mathrm{d} s+ \\
& -2 \int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left(f_{n}\left(u_{n}(s)\right), u(s)\right)+\frac{\mathbf{1}}{2}\left(\mathrm{G}_{n}\left(u_{n}(s)\right), \mathrm{G}(v(s))\right)_{*}\right\} \mathrm{d} s
\end{aligned}
$$

(14) implies:

$$
\begin{aligned}
\left|u_{n}(\mathrm{~T})\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2} & -\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})} \leq 2 \int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left(f_{n}\left(u_{n}(s), \mathrm{J}_{n}\left(u_{n}(s)\right)\right)\right\} \mathrm{d} s+\right. \\
& +\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\left\|\mathrm{G}_{n}\left(u_{n}(s)\right)\right\|_{*}^{2}\right\} \mathrm{d} s
\end{aligned}
$$

hence as $n \rightarrow \infty$
(20)

$$
\begin{aligned}
\underline{\lim } \mathrm{X}_{n} & \geq|u(\mathrm{~T})|_{\mathrm{L}^{2}(\mathrm{H})}^{2}-\left|u_{0}\right|_{\mathrm{L}^{2}(\mathrm{H})}^{2}+ \\
& -2 \int_{0}^{\mathrm{T}} \mathrm{E}\{\langle f(v(s)), u(s)-v(s)|\} \mathrm{d} s+ \\
& -\int_{0}^{\mathrm{T}} \mathrm{E}\left\{(\mathrm{G}(v(s)), \psi(s)-\mathrm{G}(v(s)))_{*}\right\} \mathrm{d} s+ \\
& -2 \int_{0}^{\mathrm{T}} \mathrm{E}\{|\chi(s), v(s)|\} \mathrm{d} s+\int_{0}^{\mathrm{T}} \mathrm{E}\left\{(\psi(s), \mathrm{G}(v(s)))_{*}\right\} \mathrm{d} s .
\end{aligned}
$$

Now from (I 9 ) it is:

$$
\begin{align*}
&|u(\mathrm{~T})|_{\mathrm{L}^{2}(\mathrm{H})}^{2}=\left|u_{0}\right|_{L^{2}\left(\mathrm{H}_{0}\right)}^{2}+2 \int_{0}^{\mathrm{T}} \mathrm{E}\{\langle\chi(s), u(s)\rangle\} \mathrm{d} s+  \tag{2I}\\
&+\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\|\psi(s)\|_{*}^{2}\right\} \mathrm{d} s
\end{align*}
$$

so that substituting in (20):

$$
\begin{align*}
& 2 \int_{0}^{\mathrm{T}} \mathrm{E}\{(\chi(s)-f(v(s)), u(s)-v(s)\rangle\} \mathrm{d} s+  \tag{22}\\
+ & \int_{0}^{\mathrm{T}} \mathrm{E}\left\{\|\psi(s)-\mathrm{G}(v(s))\|_{*}^{2}\right\} \mathrm{d} s \leq \lim \mathrm{X}_{n} \leq o .
\end{align*}
$$

This latter inequality gives:

$$
\int_{0}^{\mathrm{T}} \mathrm{E}\{|\chi(s)-f(v(s)), u(s)-v(s)|\} \mathrm{d} s \leq 0
$$

and the hemicontinuity yields, by a standard argument:

$$
\chi=f(u)
$$

On the other hand, from this, putting $v=u$ in (22) it is:

$$
\int_{0}^{\mathrm{T}} \mathrm{E}\left\{\|\psi(s)-\mathrm{G}(u(s))\|_{*}^{2}\right\} \mathrm{d} s \leq 0
$$

that is $\psi=G(u)$.

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