# ON A STOCHASTIC 2D CAHN-HILLIARD-NAVIER-STOKES SYSTEM DRIVEN BY JUMP NOISE 

G. DEUGOUÉ AND T. TACHIM MEDJO*


#### Abstract

We investigate a stochastic 2D Cahn-Hilliard-Navier-Stokes system with a multiplicative noise of Lévy type. The model consists of the Navier-Stokes equations for the velocity, coupled with a Cahn-Hilliard system for the order (phase) parameter. We prove that the system has a unique global strong solution and we derive some a priori estimates for the solution.


## 1. Introduction

It is well accepted that the incompressible Navier-Stokes equation governs the motions of single-phase fluids such as air or water. On the other hand, we are faced with the difficult problem of understanding the motion of binary fluid mixtures, that is fluids composed by either two phases of the same chemical species or phases of different composition. Diffuse interface models are well-known tools to describe the dynamics of complex (e.g., binary) fluids, $[25,26]$. For instance, this approach is used in [5] to describe cavitation phenomena in a flowing liquid. The model consists of the NSE equation coupled with the phase-field system, [14, 25, 26, 27]. In the isothermal compressible case, the existence of a global weak solution is proved in [24]. In the incompressible isothermal case, neglecting chemical reactions and other forces, the model reduces to an evolution system which governs the fluid velocity $v$ and the order parameter $\phi$. This system can be written as a NSE equation coupled with a convective Allen-Cahn equation, [25]. The associated initial and boundary value problem was studied in [25] in which the authors proved that the system generated a strongly continuous semigroup on a suitable phase space which possesses a global attractor. When the two fluids have the same constant density, the temperature differences are negligible and the diffuse interface between the two phases has a small but non-zero thickness, a well-known model is the so-called "Model H" (cf. [28, 30]). This is a system of equations where an incompressible Navier-Stokes equation for the (mean) velocity $v$ is coupled with a convective Cahn-Hilliard equation for the order parameter $\phi$, which represents the relative concentration of one of the fluids.

The purpose of this article is to study a stochastic 2D Cahn-Hilliard-NavierStokes equations (CH-NSE) driven by a non-Gaussian Levy noise. We recall that introducing a random term in a fluid model such as the Navier-Stokes system is now

[^0]a well accepted tool to model the influence of internal, external and environmental noise. Adding a noise term in a fluid model can also be used to describe systems that are too complex to be described deterministically, e.g. a flow of a chemical substance in a river subjected by wind and rain, an airflow around an airplane wing perturbed by the random state of the atmosphere and weather, a laser beam subjected to turbulent movement of the atmosphere, spread of an epidemic in some regions and the spatial spread of infectious diseases, $[4,11,12,23,31,32,33,34$, 41, 42, 43].

There are few notable works available on the stochastic CH-NSE driven by Gaussian noise. In [19], the authors considered the stochastic 3D globally modified Cahn-Hilliard-Navier-Stokes equations with multiplicative Gaussian noise. They proved the existence and uniqueness of strong solution (in the sense of partial differential equations and stochastic analysis). Moreover, they studied the asymptotic behavior of the unique solution and obtained the existence of a probabilistic weak solution for the stochastic 3D Cahn-Hilliard-Navier-Stokes equations. In [18], they also considered the asymptotic stability of the unique strong solution for the 3D globally modified Cahn-Hilliard-Navier-Stokes equations. The second author of the paper has proved the existence and uniqueness of the probabilistic strong solution for the stochastic 2D CH-NSE with multiplicative noise, [35].

In recent years, introducing a jump-type noises as Lévy-type or Poisson-type perturbations has become extremely popular for modeling natural phenomena, because these noises are very nice choice to reproduce the performance of some natural phenomena in real world models, such as some large moves and unpredictable events. There is a large amount of literature on the existence and uniqueness solutions for stochastic partial differential equations driven by jump-type noises. We refer the reader to $[12,20,21,22,36,37,38,39,40,44,45]$. However, the existing results in the literature do not cover the situation considered in this paper.

The aim of this article is to study a class of stochastic coupled CH-NSE driven by jump noise of Lévy type. To the best of our knowledge, this is the first work dealing with the stochastic version of the CH-NSE driven by jump noise. The model includes an abstract and general form of random external forces depending eventually on the velocity $v$ of the fluid and the order parameter $\phi$. We prove the existence and uniqueness of strong solutions. The proof of the existence of solution is based on a Galerkin scheme similar to that of $[29,11]$ in the case of the 2D Navier-Stokes and the 3D Lagrangian averaged Navier-Stokes equations. Let us note that the coupling between the Navier-Stokes and the Cahn-Hilliard systems introduces in the system a highly nonlinear coupling term that makes the analysis of the problems studied in this article more involved.

The article is divided as follows. In the next section we present the stochastic Cahn-Hilliard-Navier-Stokes model and its mathematical setting. We also give most of the notations and necessary preliminary used throughout this work. The main results appear in the third section, where we use a Galerkin approximation to prove the existence of strong solution. In the fourth section, we prove the pathwise uniqueness and the convergence of the whole Galerkin approximate solution.

## 2. The stochastic CH-NSE and its mathematical setting

2.1. Governing equations. We assume that the domain $\mathcal{M}$ of the fluid is a bounded domain in $\Re^{2}$. Then, we consider the system

$$
\left\{\begin{array}{l}
d v+\left[-\nu_{1} \Delta v+(v \cdot \nabla) v+\nabla p-\mathcal{K} \mu \nabla \phi\right] d t=g_{1}(t, v, \phi) d t  \tag{2.1}\\
\quad+\int_{Z} \sigma(t, v, \phi, z) \tilde{\eta}(d t, d z) \text { in }(0, T) \times \mathcal{M} \\
\operatorname{div} v=0 \text { in }(0, T) \times \mathcal{M} \\
\frac{\partial \phi}{\partial t}+v \cdot \nabla \phi-\nu_{3} \Delta \mu=0 \text { in }(0, T) \times \mathcal{M} \\
\mu=-\nu_{2} \Delta \phi+\alpha f(\phi) \text { in }(0, T) \times \mathcal{M}
\end{array}\right.
$$

In (2.1), the unknown functions are the velocity $v=\left(v_{1}, v_{2}\right)$ of the fluid, its pressure $p$ and the order (phase) parameter $\phi$.

The terms $g_{1}(t, v, \phi)$ and $\int_{Z} \sigma(t, v, \phi, z) \tilde{\eta}(d t, d z)$ respectively represent the deterministic and the random external forces that eventually depend on $(v, \phi)$, and $\tilde{\eta}$ is a compensated Poisson measure on a measurable space $(Z, \mathcal{Z})$ endowed with a fixed $\Sigma$-finite measure $\nu$. Precise assumption on the data are given below. The model (2.1) describes the motion of a binary fluid exited by random forces.

The quantity $\mu$ is the variational derivative of the following free energy functional

$$
\begin{equation*}
\mathcal{F}_{p}(\phi)=\int_{\mathcal{M}}\left(\frac{\nu_{2}}{2}|\nabla \phi|^{2}+\alpha F(\phi)\right) d s \tag{2.2}
\end{equation*}
$$

where, e.g., $F(r)=\int_{0}^{r} f(\zeta) d \zeta$. Here, the constants $\nu_{1}>0, \nu_{3}>0$ and $\mathcal{K}>$ 0 correspond to the kinematic viscosity of the fluid, the mobility constant and the capillarity (stress) coefficient respectively. Here $\nu_{2}, \alpha>0$ are two physical parameters describing the interaction between the two phases. In particular, $\nu_{2}$ is related with the thickness of the interface separating the two fluids.

A typical example of potential $F$ is that of logarithmic type.
However, this potential is often replaced by a polynomial approximation of the type $F(r)=\gamma_{1} r^{4}-\gamma_{2} r^{2}, \gamma_{1}, \gamma_{2}$ being positive constants. As noted in [25], (2.1) $)_{1}$ can be replaced by

$$
\begin{align*}
& d v(t)+\left[-\nu_{1} \Delta v+(v \cdot \nabla) v+\nabla \tilde{p}\right] d t=\int_{Z} \sigma(t, v, \phi, z) \tilde{\eta}(d t, d z)  \tag{2.3}\\
& \quad+\left[-\mathcal{K} \operatorname{div}(\nabla \phi \otimes \nabla \phi)+g_{1}(t, v, \phi)\right] d t
\end{align*}
$$

where $\tilde{p}=p-\mathcal{K}\left(\frac{\nu_{2}}{2}|\nabla \phi|^{2}+\alpha F(\phi)\right)$, since $\mathcal{K} \mu \nabla \phi=\nabla\left(\mathcal{K}\left(\frac{\nu_{2}}{2}|\nabla \phi|^{2}+\alpha F(\phi)\right)\right)-$ $\mathcal{K} \operatorname{div}(\nabla \phi \otimes \nabla \phi)$. The stress tensor $\nabla \phi \otimes \nabla \phi$ is considered the main contribution modeling capillary forces due to surface tension at the interface between the two phases of the fluid.

Regarding the boundary conditions for these models, we assume that the boundary conditions for $\phi$ are the natural no-flux condition

$$
\begin{equation*}
\partial_{\eta} \phi=\partial_{\eta} \Delta \phi=0, \quad \text { on } \partial \mathcal{M} \times(0, \infty) \tag{2.4}
\end{equation*}
$$

where $\partial \mathcal{M}$ is the boundary of $\mathcal{M}$ and $\eta$ is the outward normal to $\partial \mathcal{M}$. These conditions ensure the mass conservation. Note that (2.4) implies that

$$
\begin{equation*}
\partial_{\eta} \mu=0, \quad \text { on } \partial \mathcal{M} \times(0, \infty) \tag{2.5}
\end{equation*}
$$

From (2.5), we deduce the conservation of the following quantity

$$
\begin{equation*}
\langle\phi(t)\rangle=\frac{1}{|\mathcal{M}|} \int_{\mathcal{M}} \phi(x, t) d x \tag{2.6}
\end{equation*}
$$

where $|\mathcal{M}|$ stands for the Lebesgue measure of $\mathcal{M}$. More precisely, we have

$$
\begin{equation*}
\langle\phi(t)\rangle=\langle\phi(0)\rangle, \quad \forall t \geq 0 \tag{2.7}
\end{equation*}
$$

Concerning the boundary condition for $v$, we assume the Dirichlet (no-slip) boundary condition

$$
\begin{equation*}
v=0, \quad \text { on } \partial \mathcal{M} \times(0, \infty) \tag{2.8}
\end{equation*}
$$

Therefore we assume that there is no relative motion at the fluid-solid interface.
The initial condition is given by

$$
\begin{equation*}
(v, \phi)(0)=\left(v_{0}, \phi_{0}\right), \text { in } \mathcal{M} \tag{2.9}
\end{equation*}
$$

2.2. Mathematical setting. We first recall from [25] a weak formulation of (2.1), (2.4), (2.8)-(2.9). Hereafter, we assume that the domain $\mathcal{M}$ is bounded with a smooth boundary $\partial \mathcal{M}$ (e.g., of class $\left.\mathcal{C}^{3}\right)$. We also assume that $f \in \mathcal{C}^{2}(\Re)$ satisfies

$$
\left\{\begin{array}{l}
\lim _{|r| \rightarrow+\infty} f^{\prime}(r)>0  \tag{2.10}\\
\left|f^{\prime}(r)\right| \leq c_{f}\left(1+|r|^{k}\right), \forall r \in \Re
\end{array}\right.
$$

where $c_{f}$ is some positive constant and $k \in[1,+\infty)$ is fixed. It follows from (2.10) that

$$
\begin{equation*}
|f(r)| \leq c_{f}\left(1+|r|^{k+1}\right), \forall r \in \Re . \tag{2.11}
\end{equation*}
$$

Note that the derivative of the typical double-well potential $f$ satisfies conditions similar to (2.10). Let us now recall from [25] the functional set up of the model (2.1), (2.4), (2.8), (2.9).

If $X$ is a real Hilbert space with inner product $(\cdot, \cdot)_{X}$, we will denote the induced norm by $|\cdot|_{X}$, while $X^{*}$ will indicate its dual. We set

$$
\mathcal{V}_{1}=\left\{u \in\left(\mathcal{C}_{c}^{\infty}(\mathcal{M})\right)^{2}: \operatorname{div} u=0, \text { in } \mathcal{M}\right\}
$$

We denote by $H_{1}$ and $V_{1}$ the closure of $\mathcal{V}_{1}$ in $\left(L^{2}(\mathcal{M})\right)^{2}$ and $\left(H_{0}^{1}(\mathcal{M})\right)^{2}$ respectively. The scalar product in $H_{1}$ is denoted by $(\cdot, \cdot)_{L^{2}}$ and the associated norm by $|\cdot|_{L^{2}}$. Moreover, the space $V_{1}$ is endowed with the scalar product

$$
((u, v))=\sum_{i=1}^{2}\left(\partial_{x_{i}} u, \partial_{x_{i}} v\right)_{L^{2}}, \quad\|u\|=((u, u))^{1 / 2}
$$

We now define the operator $A_{0}$ by

$$
A_{0} u=-\mathcal{P}_{1} \Delta u, \forall u \in D\left(A_{0}\right)=\left(H^{2}(\mathcal{M})\right)^{2} \cap V_{1}
$$

where $\mathcal{P}_{1}$ is the Leray-Helmotz projector in $\left(L^{2}(\mathcal{M})\right)^{2}$ onto $H_{1}$. Then, $A_{0}$ is a self-adjoint positive unbounded operator in $H_{1}$ which is associated with the scalar product defined above. Furthermore, $A_{0}^{-1}$ is a compact linear operator on $H_{1}$ and $\left|A_{0} \cdot\right|_{L^{2}}$ is a norm on $D\left(A_{0}\right)$ that is equivalent to the $H^{2}$-norm.

We introduce the linear nonnegative unbounded operator on $L^{2}(\mathcal{M})$

$$
\begin{equation*}
A_{1} \phi=-\Delta \phi, \forall \phi \in D\left(A_{1}\right)=\left\{\phi \in H^{2}(\mathcal{M}), \partial_{\eta} \phi=0, \quad \text { on } \partial \mathcal{M}\right\} \tag{2.12}
\end{equation*}
$$

and we endow $D\left(A_{1}\right)$ with the norm $\left|A_{1} \cdot\right|_{L^{2}}+|\langle\cdot\rangle|_{L^{2}}$, which is equivalent to the $H^{2}$-norm. We also define the linear positive unbounded operator on the Hilbert space $L_{0}^{2}(\mathcal{M})$ of the $L^{2}$ - functions with null mean

$$
\begin{equation*}
B_{n} \phi=-\Delta \phi, \forall \phi \in D\left(B_{n}\right)=D\left(A_{1}\right) \cap L_{0}^{2}(\mathcal{M}) . \tag{2.13}
\end{equation*}
$$

Note that $B_{n}^{-1}$ is a compact linear operator on $L_{0}^{2}(\mathcal{M})$. More generally, we can define $B_{n}^{s}$, for any $s \in \Re$, noting that $\left|B_{n}^{s / 2} \cdot\right|_{L^{2}}, s>0$, is an equivalent norm to the canonical $H^{s}-$ norm on $D\left(B_{n}^{s / 2}\right) \subset H^{s}(\mathcal{M}) \cap L_{0}^{2}(\mathcal{M})$. Also note that $A_{1}=B_{n}$ on $D\left(B_{n}\right)$. If $\phi$ is such that $\phi-\langle\phi\rangle \in D\left(B_{n}^{s / 2}\right)$, we have that $\left|B_{n}^{s / 2}(\phi-\langle\phi\rangle)\right|_{L^{2}}+|\langle\phi\rangle|_{L^{2}}$ is equivalent to the $H^{s}$-norm. Moreover, we set $H^{-s}(\mathcal{M})=\left(H^{s}(\mathcal{M})\right)^{*}$, whenever $s<0$.

Finally we set

$$
\begin{equation*}
H_{2}=D\left(B_{n}^{0}\right)=L_{0}^{2}(\mathcal{M}), V_{2}=D\left(B_{n}^{1 / 2}\right) . \tag{2.14}
\end{equation*}
$$

The norms in $H_{2}$ and $V_{2}$ are denoted respectively by $|\cdot|_{L^{2}}$ and $\|\cdot\|$, where $\|\psi\|=$ $\left|B_{n}^{1 / 2} \psi\right|_{L^{2}}$.

We introduce the bilinear operators $B_{0}, B_{1}$ (and their associated trilinear forms $\left.b_{0}, b_{1}\right)$ as well as the coupling mapping $R_{0}$, which are defined from $D\left(A_{0}\right) \times D\left(A_{0}\right)$ into $H_{1}, D\left(A_{0}\right) \times D\left(A_{1}\right)$ into $L^{2}(\mathcal{M})$, and $\left(L^{2}(\mathcal{M})\right)^{2} \times\left(D\left(A_{1}\right) \cap H^{3}(\mathcal{M})\right)$ into $H_{1}$, respectively. More precisely, we set

$$
\begin{aligned}
& \left\langle B_{0}(u, v), w\right\rangle=\int_{\mathcal{M}}[(u \cdot \nabla) v] \cdot w d x=b_{0}(u, v, w), \forall u, v, w \in D\left(A_{0}\right), \\
& \left\langle B_{1}(u, \phi), \rho\right\rangle=\int_{\mathcal{M}}[(u \cdot \nabla) \phi] \rho d x=b_{1}(u, \phi, \rho), \forall u \in D\left(A_{0}\right), \phi, \rho \in D\left(A_{1}\right), \\
& \left\langle R_{0}(\mu, \phi), w\right\rangle=\int_{\mathcal{M}} \mu[\nabla \phi \cdot w] d x=b_{1}(w, \phi, \mu), \\
& \quad \forall w \in D\left(A_{0}\right), \phi \in D\left(A_{1}\right) \cap H^{3}(\mathcal{M}), \mu \in L^{2}(\mathcal{M}) .
\end{aligned}
$$

Note that

$$
R_{0}(\mu, \phi)=\mathcal{P} \mu \nabla \phi .
$$

We recall from [25] (see also [26, 27]) that $B_{0}, B_{1}$ and $R_{0}$ satisfy the following estimates

$$
\begin{align*}
& \left|B_{0}(u, v)\right|_{V_{1}^{*}} \leq c|u|_{L^{2}}^{1 / 2}\|u\|^{1 / 2}\|v\|, \forall u, v \in V_{1}, \\
& \left|B_{0}(u, v)\right|_{L^{2}} \leq c|u|_{L^{2}}^{1 /}\|u\|^{1 / 2}\|v\|^{1 / 2}\left|A_{0} v\right|_{L^{2}}^{1 / 2}, \forall u \in V_{1}, v \in D\left(A_{0}\right),  \tag{2.15}\\
& \left|B_{1}(u, \phi)\right|_{V_{2}^{*}} \leq c|u|_{L^{2}}^{1 / 2}\|u\|^{1 / 2}\|\phi\|, \forall u \in V_{1}, \phi \in V_{2},  \tag{2.16}\\
& \left|B_{1}(u, \phi)\right|_{L^{2}} \leq c|u|_{L^{2}}^{1 / 2}\|u\|^{1 / 2}\|\phi\|^{1 / 2}\left|A_{1} \phi\right|_{L^{2}}^{1 / 2}, \forall u \in V_{1}, \phi \in D\left(A_{1}\right), \\
& \quad\left|R_{0}\left(A_{1} \phi, \rho\right)\right|_{V_{1}^{*}} \leq c\left|A_{1} \phi\right|_{L^{2}}^{1 / 2}|\phi|_{H^{3}}^{1 / 2}\|\rho\|, \forall \phi \in D\left(A_{1}\right), \rho \in V_{2}, \\
& \quad\left|R_{0}\left(A_{1} \phi, \rho\right)\right|_{L^{2}} \leq c\|\rho\| \|^{1 / 2}\left|A_{1} \rho\right|_{L^{2}}^{1 / 2}\left|A_{1} \phi\right|_{L^{2}}^{1 / 2}|\phi|_{H^{3}}^{1 / 2},  \tag{2.17}\\
& \quad \forall \phi \in D\left(A_{1}\right), \rho \in D\left(A_{1}^{3 / 2}\right) .
\end{align*}
$$

We recall that (due to the mass conservation) we have

$$
\begin{equation*}
\langle\phi(t)\rangle=\langle\phi(0)\rangle=M_{0}, \forall t>0 . \tag{2.18}
\end{equation*}
$$

Thus, up to a shift of the order parameter field, we can always assume that the mean of $\phi$ is zero a the initial time and, therefore it will remain zero for all positive times. Hereafter, we assume that

$$
\begin{equation*}
\langle\phi(t)\rangle=\langle\phi(0)\rangle=0, \forall t>0 \tag{2.19}
\end{equation*}
$$

We set

$$
\begin{equation*}
\mathcal{H}=H_{1} \times D\left(B_{n}^{1 / 2}\right) \tag{2.20}
\end{equation*}
$$

The space $\mathcal{H}$ is a complete metric space with respect to the norm

$$
\begin{equation*}
|(v, \phi)|_{\mathcal{H}}^{2}=\mathcal{K}^{-1}|v|_{L^{2}}^{2}+\nu_{2}|\nabla \phi|_{L^{2}}^{2} . \tag{2.21}
\end{equation*}
$$

We define the Hilbert space $\mathcal{U}$ by

$$
\begin{equation*}
\mathcal{U}=V_{1} \times D\left(B_{n}\right) \tag{2.22}
\end{equation*}
$$

endowed with the scalar product whose associated norm is

$$
\begin{equation*}
\|(v, \phi)\|_{\mathcal{U}}^{2}=\|v\|^{2}+\left|B_{n} \phi\right|_{L^{2}}^{2} \tag{2.23}
\end{equation*}
$$

We will also denote by $c$ a generic positive constant that depends on the domain $\mathcal{M}$. To simplify the notations, we set (without loss of generality) $\nu_{1}=\nu_{2}=\nu_{3}=$ $\alpha=\mathcal{K}=1$.

Using the notations above, we rewrite (2.1), (2.4), (2.8)-(2.6) as

$$
\left\{\begin{array}{l}
d v+\left[A_{0} v+B_{0}(v, v)-R_{0}\left(A_{1} \phi, \phi\right)\right] d t=g_{1}(t, v, \phi) d t  \tag{2.24}\\
\quad+\int_{Z} \sigma(t, v, \phi, z) \tilde{\eta}(d t, d z) \\
\frac{d \phi}{d t}+A_{1} \mu+B_{1}(v, \phi)=0 \\
\mu=A_{1} \phi+f(\phi) \\
(v, \phi)(0)=\left(v_{0}, \phi_{0}\right)
\end{array}\right.
$$

Remark 2.1. In the weak formulation (2.24), the term $\mu \nabla \phi$ is replaced by $A_{1} \nabla \phi$. This is justified since $f^{\prime}(\phi) \nabla \phi$ is the gradient $F(\phi)$ and can be incorporated into the pressure gradient, see [25] for details. For the sake of convenience, as in [25] we will replace $\mu$ in $(2.24)_{3}$ by $\bar{\mu}=\mu-\langle\mu\rangle$, that is $\bar{\mu}=A_{1} \phi+f(\phi)-$ $\langle f(\phi)\rangle$, a.e., in $\mathcal{M} \times(0, T)$. Obviously we have $\langle\bar{\mu}(t)\rangle=0, \forall t>0$.

Notations. We first recall from [29, 11] some notations and stochastic preliminaries.

Hereafter, by $\aleph$ we denote the set of nonnegative integers, i.e. $\aleph=\{0,1,2, \cdots\}$ and by $\bar{\aleph}$ we denote the set $\aleph \cup\{+\infty\}$. Whenever we speak about $\aleph ~($ or $\bar{\aleph})$-valued measurable functions we implicitly assume that the set is equipped with the trivial $\Sigma$-field $2^{\aleph}$ (or $2^{\bar{\aleph}}$ ). By $\Re_{+}$we will denote the interval $[0, \infty)$ and by $\Re_{*}$ the set $\Re \backslash$ $\{0\}$. If $X$ is a topological space, then by $\mathcal{B}(X)$ we will denote the Borel $\Sigma$-field on $X$. By $\lambda_{d}$ we will denote the Lebesgue measure on $\left(\Re^{d}, \mathcal{B}\left(\Re^{d}\right)\right)$, by $\lambda$ the Lebesgue measure on ( $\Re, \mathcal{B}(\Re)$ ).

If $(S, \mathfrak{S})$ is a measurable space then by $M(S)$ we denote the set of all real valued measures on $(S, \mathfrak{S})$, and by $M(S)$ the $\Sigma$-field on $M(S)$ generated by the functions $i_{B}: M(S) \ni \varsigma \mapsto \varsigma(B) \in \Re, B \in S$. By $M_{+}(S)$ we denote the set of all nonnegative measures on S , and by $M(S)$ the $\Sigma$-field on $M_{+}(S)$ generated by the functions
$i_{B}: M_{+}(S) \ni \varsigma \mapsto \varsigma(B) \in \Re_{+}, B \in S$. Finally, by $M_{I}(S)$ we denote the family of all $\bar{\aleph}$ valued measures on $(S, \mathfrak{S})$, and by $M_{I}(S)$ the $\Sigma$-field on $M_{I}(S)$ generated by functions $i_{B}: M(S) \ni \varsigma \mapsto \varsigma(B) \in \bar{\aleph}, B \in S$. If $(S, \mathfrak{S})$ is a measurable space then we will denote by $\mathfrak{S} \otimes \mathcal{B}\left(\Re_{+}\right)$the product $\Sigma$-field on $S \times \Re_{+}$and by $\nu \otimes \lambda$ the product measure of $\nu$ and the Lebesgue measure $\lambda$.

Preliminaries. As mentioned earlier we will study a stochastic model for a CHNSE excited by random forces. We first describe the forces acting on the fluids. Let $(Z, \mathcal{Z})$ be a separable metric space and let $\nu$ be a $\Sigma$-finite positive measure on it. Suppose that $\mathfrak{P}=(\Omega, \mathfrak{F}, \mathbb{F}, \mathbb{P})$ is a filtered probability space, where $\mathbb{F}=(\mathfrak{F})_{t \geq 0}$ is a filtration satisfying the usual conditions, and $\eta: \Omega \times \mathcal{B}\left(\Re_{+}\right) \times Z \rightarrow \bar{\aleph}$ is a time homogeneous Poisson random measure, with intensity measure $\nu$, defined over the filtered probability space $\mathfrak{P}$. A time homogeneous Poisson random measure defined over $\mathfrak{P}$ is given in the following definition.

Definition 2.2. Let $Z$ be a metric space and $\mathcal{Z}$ its Borel $\Sigma$-algebra, $\nu$ a positive $\Sigma$-finite measure on $(Z, \mathcal{Z})$. A Poisson random measure, with intensity measure $\nu$ defined on $(Z, \mathbb{Z})$ over $\mathfrak{P}$ is a measurable map $\eta:(\Omega, \mathfrak{F}) \rightarrow\left(M_{I}\left(Z \times \Re_{+}\right), M_{I}(Z \times\right.$ $\Re+)$ ) satisfying the following conditions:
(i) for all $B \in \mathcal{B}\left(Z \otimes \Re_{+}\right), \eta(B): \Omega \rightarrow \bar{\aleph}$ is a Poisson random measure with parameter $\mathbb{E}[\eta(B)]$;
(ii) $\eta$ is independently scattered, i.e., if the sets $B_{j} \in \mathcal{B}\left(Z \otimes \Re_{+}\right), j=1, \ldots, n$, are disjoint then the random variables $\eta\left(B_{j}\right), j=1, \cdots, n$, are independent;
(iii) for all $U \in \mathcal{Z}$ and $I \in \mathcal{B}\left(\Re_{+}\right)$

$$
\mathbb{E}[\eta(U \times I)]=\lambda(I) \nu(U)
$$

(iv) for all $U \in \mathcal{Z}$ the $\bar{\aleph}-$ valued process $(N(U, t))_{t \geq 0}$ defined by $N(U, t):=\eta(U \times$ $(0, t]), t \geq 0$, is $\mathbb{F}$-adapted and its increments are independent of the past, i.e., if $t>s \geq 0$, then the random variable $N(U, t)-N(U, s)=\eta(U \times(s, t])$ is independent of $\mathfrak{F}_{s}$.

We will denote by $\tilde{\eta}$ the compensated Poisson random measure defined by

$$
\tilde{\eta}:=\eta-\gamma
$$

where the compensator $\gamma: \mathcal{B}\left(\mathcal{Z} \times \Re_{+}\right) \rightarrow \Re_{+}$is defined by

$$
\gamma(A \times I)=\lambda(I) \nu(A), I \in \mathcal{B}\left(\Re_{+}\right), A \in \mathcal{Z}
$$

As noted in [29], while items (i) and (ii) are the classical definition, see for e.g. Definition 6.1 in [39], of a Poisson Random measure $\eta$, the remaining items implicitly indicate that $\eta$ is associated to a certain Lévy process $\tilde{L}$; see, for instance [[39], Proposition 4.16].

Let $\mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$ be the class of all progressively measurable processes $\xi: \Re_{+} \times Z \times \Omega \rightarrow H_{1}$ satisfying the condition

$$
\begin{equation*}
\mathbb{E} \int_{0}^{T} \int_{Z}|\xi(r, z)|_{L^{2}}^{2} \nu(d z) d r<\infty, \forall T>0 \tag{2.25}
\end{equation*}
$$

If $T>0$, the class of all progressively measurable processes $\xi:[0, T] \times Z \times$ $\Omega \rightarrow H_{1}$ satisfying the condition (2.25) just for this one $T$, will be denoted by
$\mathcal{M}^{2}\left(0, T, L^{2}\left(Z, \nu, H_{1}\right)\right)$. Also, let $\mathcal{M}_{\text {step }}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$ be the space of all processes $\xi \in \mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$ such that

$$
\begin{equation*}
\xi(r)=\sum_{j=1}^{n} 1_{\left(t_{j-1}, t_{j}\right]}(r) \xi_{j}, 0 \leq r \tag{2.26}
\end{equation*}
$$

where $\left\{0=t_{0}<t_{1}<\ldots<t_{n}<\infty\right\}$ is a partition of $[0, \infty)$, and for all $j, \xi_{j}$ is an $\mathfrak{F}_{t_{j-1}}-$ measurable random variable.

For any $\xi \in \mathcal{M}_{\text {step }}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$, we set

$$
\tilde{I}(\xi)=\sum_{j=1}^{n} \int_{Z} \xi_{j}(z) \tilde{\eta}\left(d z,\left(t_{j-1}, t_{j}\right]\right)
$$

This is basically the definition of stochastic integral of a random step process $\xi$ with respect to the compound random Poisson measure $\tilde{\eta}$. The extension of this integral on $\mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$ is possible thanks to the following result which is taken from [39].

Theorem 2.3. There exists a unique bounded linear operator

$$
I: \mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu ; H_{1}\right)\right) \rightarrow L^{2}\left(\Omega, \mathfrak{F} ; H_{1}\right)
$$

such that for $\xi \in \mathcal{M}_{\text {step }}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$ we have $\tilde{I}(\xi)=I(\xi)$. In particular, there exists a constant $C=C\left(H_{1}\right)$ such that for any $\xi \in \mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$,

$$
\mathbb{E}\left|\int_{0}^{t} \int_{Z} \xi(r, z) \tilde{\eta}(d z, d r)\right|_{L^{2}}^{2} \leq C \mathbb{E} \int_{0}^{t} \int_{Z}|\xi(r, z)|_{L^{2}}^{2} \eta(d z) d r, t>0
$$

Moreover, for each $\xi \in \mathcal{M}^{2}\left(\Re_{+}, L^{2}\left(Z, \nu, H_{1}\right)\right)$, the process $I\left(1_{[0, t]} \xi\right), t \geq 0$, is an $H_{1}$-valued càdlàg martingale. The process $1_{[0, t]} \xi$ is defined by $\left[1_{[0, t]} \xi\right](r, z, \omega):=$ $1_{[0, t]}(r)(r, z, \omega), t \geq 0, r \in \Re_{+}, z \in Z$ and $\omega \in \Omega$.

As usual we will write

$$
\int_{0}^{t} \int_{Z} \xi(r, z) \tilde{\eta}(d z, d r):=I(\xi)(t), t \geq 0
$$

If $T>0$, we denote by $\mathbb{D}\left(0, T ; H_{1}\right)$ the space of all càdlàg paths from $[0, T]$ into $H_{1}$.

Now we introduce the main set of hypotheses used in this article. As in [29, 11], we suppose that we are given a function $\sigma$ satisfying the following set of constraints:

Condition 1. There exist nonnegative constants $l_{0}, l_{1}, l_{2}$ such that, for any $t \in[0, T]$ and all $\left(v_{1}, \phi_{1}\right),\left(v_{2}, \phi_{2}\right) \in \mathcal{H}$, we have

$$
\begin{align*}
& \left|\sigma\left(t, v_{1}, \phi_{1}\right)\right|_{L^{2}\left(Z, \nu ; H_{1}\right)}^{p} \leq l_{0}+l_{1}\left|\left(v_{1}, \phi_{1}\right)\right|_{\mathcal{H}}^{p} ; \quad \text { for any } p \geq 2, \\
& \left|\sigma\left(t, v_{1}, \phi_{1}\right)-\sigma\left(t, v_{2}, \phi_{2}\right)\right|_{L^{2}\left(Z, \nu ; H_{1}\right)}^{2} \leq l_{2}\left|\left(v_{1}, \phi_{1}\right)-\left(v_{2}, \phi_{2}\right)\right|_{\mathcal{H}}^{2} . \tag{2.27}
\end{align*}
$$

We assume that the external forcing $g_{1}$ is a measurable Lipschitz and sublinear mappings from $\Omega \times(0, T) \times H_{1}$ into $V_{1}^{*}$. More precisely, for all $\left(v_{1}, \phi_{1}\right),\left(v_{2}, \phi_{2}\right) \in V_{1}$,
$g_{1}\left(\cdot, v_{1}, \phi_{1}\right)$ is $\mathfrak{F}_{t}-$ progressively measurable, and $d \mathbb{P} \times d t-$ a.e. in $\Omega \times(0, T)$

$$
\begin{align*}
& \left\|g_{1}\left(t, v_{1}, \phi_{1}\right)-g_{1}\left(t, v_{2}, \phi_{2}\right)\right\|_{V_{1}^{*}} \leq L_{1}\left|v_{1}-v_{2}\right|_{L^{2}},  \tag{2.28}\\
& g_{1}(t, 0,0) \in M_{\widetilde{\mathcal{F}}_{t}}^{2}\left(0, T ; V_{1}^{*}\right) .
\end{align*}
$$

Finally, we assume that

$$
\begin{equation*}
\left(v_{0}, \phi_{0}\right) \in L^{2}\left(\Omega, \mathfrak{F}_{0}, \mathbb{P} ; \mathcal{H}\right) \tag{2.29}
\end{equation*}
$$

Hereafter, for any $(w, \psi) \in \mathcal{H}$, we set

$$
\begin{equation*}
\mathcal{E}(w, \psi)=|w|_{L^{2}}^{2}+\|\psi\|^{2}+2\langle F(\psi), 1\rangle+c_{1} \tag{2.30}
\end{equation*}
$$

where $c_{1}>0$ is a constant large enough and independent on $(w, \psi)$ such that $\mathcal{E}(w, \psi)$ is nonnegative (note that $F$ is bounded from below).

We can check that (see [25]) there exists a monotone non-decreasing function $Q_{0}$ (independent on time and the initial condition) such that

$$
\begin{equation*}
|(w, \psi)|_{\mathcal{H}}^{2} \leq \mathcal{E}(w, \psi) \leq Q_{0}\left(|(w, \psi)|_{\mathcal{H}}^{2}\right), \quad \forall(w, \psi) \in \mathcal{H} \tag{2.31}
\end{equation*}
$$

Definition 2.4. Let $(Z, \mathcal{Z})$ be a separable metric space on which is defined a $\Sigma$-finite measure $\nu$ and $\left(v_{0}, \phi_{0}\right) \in L^{2}\left(\Omega, \mathcal{F}_{0}, \mathbb{P} ; \mathcal{H}\right)$. A strong solution to the problem (2.24) is a stochastic process $(v, \phi)$ such that
(1) $(v, \phi)=\{(v, \phi)(t), t \geq 0\}$ is a $\mathbb{F}$-progressively measurable process such that

$$
\mathbb{E} \sup _{s \in[0, T]} \mathcal{E}(v(s), \phi(s))+\mathbb{E} \int_{0}^{T}\|(v, \phi)(s)\|_{\mathcal{U}}^{2} d t<\infty
$$

(2) the following holds

$$
\begin{align*}
& (v(t), w)=\left(v_{0}, w\right)+\int_{0}^{t} \int_{Z}\langle\sigma(s, v, \phi, z), w\rangle \tilde{\eta}(d z, d s) \\
& \quad-\int_{0}^{t}\left\langle A_{0} v+B_{0}(v, v)-R_{0}\left(A_{1} \phi, \phi\right)-g_{1}(s, v, \phi), w\right\rangle d s, \forall w \in V_{1}  \tag{2.32}\\
& (\phi(t), \psi)=\left(\phi_{0}, \psi\right)-\int_{0}^{t}\left\langle A_{1} \mu+B_{1}(v, \phi), \psi\right\rangle d s, \quad \forall \psi \in V_{2} \\
& \mu=A_{1} \phi+f(\phi)
\end{align*}
$$

for almost all $t \in[0, T]$ and $\mathbb{P}$-almost surely.
In the deterministic case, the weak formulation of (2.24) was proposed and studied in $[8,6,7,26,25]$ (see also $[2,1,15]$ ), where the existence and uniqueness results for weak and strong solutions were proved in the deterministic case.

Before we prove this result let us recall an important statement which is borrowed from [17].

Lemma 2.5. Let $X, Y, I$ and $\varphi$ be non-negative processes and $Z_{1}$ be a non-negative integrable random variable. Assume that $I$ is non-decreasing and that there exist non-negative constants $C, \alpha_{1}, \beta, \gamma_{1}, \delta_{1}$ and $T$ satisfying first

$$
\int_{0}^{T} \varphi(s) d s \leq C, \quad \text { a.s., } \quad 2 \beta_{1} e^{C} \leq 1,2 \delta_{1} e^{C} \leq \alpha_{1}
$$

and secondly for all $t \in[0, T]$ there exists a constant $C_{1}>0$ such that

$$
\begin{aligned}
& X(t)+\alpha_{1} Y(t) \leq Z_{1}+\int_{0}^{t} \varphi(r) X(r) d r+I(t), \text { a.s. } \\
& \mathbb{E} I(t) \leq \beta \mathbb{E} X(t)+\gamma_{1} \int_{0}^{t} \mathbb{E} X(s) d s+\delta_{1} \mathbb{E} Y(t)+C_{1}
\end{aligned}
$$

If $X \in L^{\infty}([0, T] \times \Omega)$, then we have

$$
\mathbb{E}\left[X(t)+\alpha_{1} Y(t)\right] \leq 2 \exp \left(C+2 t \gamma_{1} e^{C}\right)\left(\mathbb{E} Z+C_{1}\right), t \in[0, T]
$$

## 3. Existence and uniqueness of solutions

In this section, we prove the existence and the pathwise uniqueness of variational solution to (2.24). We first prove the following energy type equality.

Proposition 3.1. If $(v, \phi)$ is a variational solution to (2.24), then $(v, \phi)$ satisfies

$$
\begin{align*}
& \mathcal{E}(v, \phi)(t)+\int_{0}^{t}\left(2\|v(s)\|^{2}+\|\mu(s)\|^{2}\right) d s=\mathcal{E}\left(v_{0}, \phi_{0}\right) \\
& \quad+2 \int_{0}^{t}\left\langle g_{1}(s, v(s), \phi(s)), v(s)\right\rangle d s+\int_{0}^{t} \int_{Z} \Upsilon(s, z) \eta(d z, d s)  \tag{3.1}\\
& \quad+2 \int_{0}^{t} \int_{Z}(v(s-), \sigma(s, v(s), \phi(s), z)) \tilde{\eta}(d z, d s)
\end{align*}
$$

where

$$
\begin{align*}
& \Upsilon(s, z)=|v(s-)+\sigma(s, v(s), \phi(s), z)|_{L^{2}}^{2}-|v(s-)|_{L^{2}}^{2}  \tag{3.2}\\
& -2(v(-s), \sigma(s, v(s), \phi(s), z))
\end{align*}
$$

Proof. Applying Itô's formula to $|v(t)|_{L^{2}}^{2}$, we derive that

$$
\begin{align*}
& \left.|v(t)|_{L^{2}}^{2}=\left|v_{0}\right|_{L^{2}}^{2}+2 \int_{0}^{t}\left(-\|v(s)\|^{2}+g_{1}(s, v(s), \phi(s)), v(s)\right\rangle\right) d s \\
& \quad+2 \int_{0}^{t}\left\langle R_{0}\left(A_{1} \phi(s), \phi(s)\right), v(s)\right\rangle d s  \tag{3.3}\\
& \quad+2 \int_{0}^{t} \int_{Z}(v(s-), \sigma(s, v(s), \phi(s), z)) \tilde{\eta}(d z, d s)+\int_{0}^{t} \int_{Z} \Upsilon(s, z) \eta(d z, d s)
\end{align*}
$$

where $\Upsilon(s, z)$ is given by (3.2).
Now multiplying $(2.24)_{2}$ with $\mu$, we obtain

$$
\begin{equation*}
\frac{d}{d t}\left(\|\phi(t)\|^{2}+2\langle F(\phi(t)), 1\rangle+c_{1}\right)+2\left\langle B_{1}(v(t), \phi(t)), \mu(t)\right\rangle=0 \tag{3.4}
\end{equation*}
$$

Integrating (3.4) and adding the result to (3.3) gives

$$
\begin{align*}
& \mathcal{E}(v, \phi)(t)=\mathcal{E}\left(v_{0}, \phi_{0}\right)-\int_{0}^{t}\left(2\|v(s)\|^{2}+\|\mu(s)\|^{2}\right) d s \\
& \quad+2 \int_{0}^{t}\left\langle g_{1}(s, v(s), \phi(s)), v(s)\right\rangle d s  \tag{3.5}\\
& \quad+2 \int_{0}^{t} \int_{Z}(v(s-), \sigma(s, v(s), \phi(s), z)) \tilde{\eta}(d z, d s)+\int_{0}^{t} \int_{Z} \Upsilon(s, z) \eta(d z, d s),
\end{align*}
$$

where $\Upsilon(s, z)$ is given by (3.2).

Note that we use the properties of $B_{0}, B_{1}$ and $R_{0}$ given in (2.15). In particular, we used the fact that ( see [25])

$$
\left\langle B_{0}(v, v), v\right\rangle=0,\left\langle R_{0}\left(A_{1} \phi, \phi\right), v\right\rangle=\left\langle B_{1}(v, \phi), \mu\right\rangle=\left\langle B_{1}(v, \phi), A_{1} \phi\right\rangle .
$$

We also use the fact that

$$
\begin{align*}
\left\langle\phi^{\prime}, \mu\right\rangle & =\frac{d}{d t}\left(\frac{1}{2}\|\phi\|^{2}+\langle F(\phi(t)), 1\rangle\right)  \tag{3.6}\\
& =\frac{d}{d t}\left(\frac{1}{2}\|\phi\|^{2}+\langle F(\phi(t)), 1\rangle+\frac{1}{2} c_{1}\right) .
\end{align*}
$$

where $c_{1}$ is the the constant that appears in (2.30).
3.1. Uniqueness of solutions. The following result implies the pathwise uniqueness of the variational solutions $(v, \phi)$ to (2.24).

Proposition 3.2. Let $\left(v_{0}^{1}, \phi_{0}^{1}\right),\left(v_{0}^{2}, \phi_{0}^{2}\right)$ be two $\mathfrak{F}_{0}$-measurable and square integrable $\mathcal{H}$-valued random variables. Let $\left(v_{1}, \phi_{1}\right),\left(v_{2}, \phi_{2}\right)$ be the variational solution to (2.24) corresponding to $\left(v_{0}^{1}, \phi_{0}^{1}\right),\left(v_{0}^{2}, \phi_{0}^{2}\right)$ respectively. Then there exists a constant $C>0$ such that

$$
\begin{equation*}
\mathbb{E} \delta(t)\left|\left(v_{1}, \phi_{1}\right)-\left(v_{2}, \phi_{2}\right)\right|_{\mathcal{H}}^{2} \leq C \mathbb{E}\left|\left(v_{0}^{1}, \phi_{0}^{1}\right)-\left(v_{0}^{2}, \phi_{0}^{2}\right)\right|_{\mathcal{H}}^{2}, \tag{3.7}
\end{equation*}
$$

for all $t \in[0, T]$, where $\delta(t)$ is defined by (3.22).
Moreover, if $\left(v_{0}^{1}, \phi_{0}^{1}\right)=\left(v_{0}^{2}, \phi_{0}^{2}\right)$ almost surely, then for any $t \in[0, T]$,

$$
\begin{equation*}
\mathbb{P}\left(\left(v_{1}, \phi_{1}\right)(t)=\left(v_{2}, \phi_{2}\right)(t)\right)=1 \tag{3.8}
\end{equation*}
$$

Proof. Let $\left(v_{1}, \phi_{1}\right),\left(v_{2} . \phi_{2}\right)$ be variational solutions to (2.24). Let $(w, \psi, \mu)=$ $\left(v_{1}, \phi_{1}, \mu_{1}\right)-\left(v_{2}, \phi_{2}, \mu_{2}\right), \bar{\mu}=\mu-\langle\mu\rangle$. Then $(w, \psi)$ satisfies

$$
\left\{\begin{array}{l}
d w+\left[A_{0} w+B_{0}\left(v_{2}, w\right)+B_{0}\left(w, v_{1}\right)\right] d t=\left[g_{1}\left(t, v_{1}, \phi_{1}\right)-g_{1}\left(t, v_{2}, \phi_{2}\right)\right] d t \\
\quad+\left[R_{0}\left(A_{1} \phi_{2}, \psi\right)+R_{0}\left(A_{1} \psi, \phi_{1}\right)\right] d t \\
\quad+\int_{Z}\left(\sigma\left(t, v_{1}, \phi_{1}, z\right)-\sigma\left(t, v_{2}, \phi_{2}, z\right)\right) \tilde{\eta}(d z, d t),  \tag{3.9}\\
\frac{d \psi}{d t}+A_{1} \bar{\mu}+B_{1}\left(v_{2}, \psi\right)+B_{1}\left(w, \phi_{1}\right)=0, \\
\mu=A_{1} \psi+f\left(\phi_{1}\right)-f\left(\phi_{2}\right), \\
(w, \psi)(0)=(0,0)
\end{array}\right.
$$

Reasoning as in the proof of Proposition (3.1) above, applying Itô's formula to $|w|_{L^{2}}^{2}$ and using (3.9) $)_{1}$, we derive that

$$
\begin{align*}
& |w|_{L^{2}}^{2}+2 \int_{0}^{t}\left(\|w\|^{2}+b_{0}\left(w, v_{1}, w\right)\right) d s \\
& \quad=2 \int_{0}^{t}\left\langle g_{1}\left(t, v_{1}, \phi_{1}\right)-g_{1}\left(t, v_{2}, \phi_{2}\right), w\right\rangle d s \\
& \quad+2 \int_{0}^{t}\left\langle R_{0}\left(A_{1} \phi_{2}, \psi\right)+R_{0}\left(A_{1} \psi, \phi_{1}\right), w\right\rangle d s  \tag{3.10}\\
& \quad+2 \int_{0}^{t} \int_{Z}\left(w(s-), \sigma\left(s, v_{1}(s), \phi_{1}(s), z\right)-\sigma\left(s, v_{2}(s), \phi_{2}(s), z\right)\right) \tilde{\eta}(d z, d s) \\
& \left.\quad+\int_{0}^{t} \int_{Z} \mid \sigma\left(s, v_{1}(s), \phi_{1}(s), z\right)-\sigma\left(s, v_{2}(s), \phi_{2}(s), z\right)\right)\left.\right|_{L^{2}} ^{2} \eta(d s, d z)
\end{align*}
$$

Now we take the duality of $(3.9)_{2}$ and (3.9) $)_{3}$ with $A_{1} \bar{\mu}-\zeta A_{1} \psi$ and $A_{1} \psi$ respectively, where $\zeta>0$ is small enough and will be selected later. Adding the resulting equality to (3.10), we derive that

$$
\begin{align*}
&|w(t)|^{2}+\|\psi(t)\|^{2}+2 \int_{0}^{t}\left(\|w\|^{2}+\zeta\left|A_{1} \psi\right|_{L^{2}}^{2}+\|\mu\|^{2}\right) d s \\
&=2 \int_{0}^{t}\left(\left\langle R_{0}\left(A_{1} \phi_{2}, \psi\right)+R_{0}\left(A_{1} \psi, \phi_{1}\right), w\right\rangle-b_{0}\left(w, v_{1}, w\right)\right) d s \\
&-2 \int_{0}^{t}\left(b_{1}\left(w, \phi, A_{1} \psi\right)+b_{1}\left(v, \psi, A_{1} \psi\right)\right) d s+\int_{0}^{t} \zeta\left\langle\bar{\mu}, A_{1} \psi\right\rangle d s \\
&+2 \int_{0}^{t}\left[\zeta\left\langle f\left(\phi_{1}\right)-f\left(\phi_{2}\right), A_{1} \psi\right\rangle-\left\langle f\left(\phi_{1}\right)-f\left(\phi_{2}\right), A_{1} \bar{\mu}\right\rangle\right] d s  \tag{3.11}\\
& \quad+2 \int_{0}^{t}\left\langle g_{1}\left(t, v_{1}, \phi_{1}\right)-g_{1}\left(t, v_{2}, \phi_{2}\right), w\right\rangle d s \\
& \quad+2 \int_{0}^{t} \int_{Z}\left(w(s-), \sigma\left(s, v_{1}(s), \phi_{1}(s), z\right)-\sigma\left(s, v_{2}(s), \phi_{2}(s), z\right)\right) \tilde{\eta}(d z, d s) \\
&\left.\quad+\int_{0}^{t} \int_{Z} \mid \sigma\left(s, v_{1}(s), \phi_{1}(s), z\right)-\sigma\left(s, v_{2}(s), \phi_{2}(s), z\right)\right)\left.\right|_{L^{2}} ^{2} \eta(d z, d s) .
\end{align*}
$$

Note that

$$
\begin{gather*}
\left|b_{0}\left(w, v_{1}, w\right)\right| \leq \frac{1}{8}\|w\|^{2}+c\left\|v_{1}\right\|^{2}|w|_{L^{2}}^{2},  \tag{3.12}\\
\left|\left\langle R_{0}\left(A_{1} \psi, \phi_{1}\right), w\right\rangle\right|=\left|b_{1}\left(w, \phi_{1}, A_{1} \psi\right)\right| \\
\left.\leq \frac{1}{8}\|w\|^{2}+\zeta\left|A_{1} \psi\right|_{L^{2}}^{2}\right)+c|w|_{L^{2}}^{2}\left\|\phi_{1}\right\|^{2}\left|A_{1} \phi_{1}\right|_{L^{2}}^{2},  \tag{3.13}\\
\left|\left\langle R_{0}\left(A_{1} \phi_{2}, \psi\right), w\right\rangle\right|=\left|b_{1}\left(w, \psi, A_{1} \phi_{2}\right)\right|  \tag{3.14}\\
\leq \frac{1}{8}\left(\|w\|^{2}+\zeta\left|A_{1} \psi_{2}\right|_{L^{2}}^{2}\right)+c\left(|w|_{L^{2}}^{2}+|\nabla \psi|_{L^{2}}^{2}\right)\left\|\phi_{2}\right\|^{2}\left|\phi_{2}\right|_{H^{2}}^{2}, \\
\zeta\left|\left\langle f\left(\phi_{1}\right)-f\left(\phi_{2}\right), A_{1} \psi\right\rangle\right|  \tag{3.15}\\
\leq \frac{\zeta}{8}\left|A_{1} \psi\right|_{L^{2}}^{2} \\
+Q_{1}\left(\left|\phi_{1}\right|_{H^{1}},\left|\phi_{2}\right|_{H^{1}}\right)\|\psi\|^{2},  \tag{3.16}\\
\left|\left\langle f\left(\phi_{1}\right)-f\left(\phi_{2}\right), A_{1} \bar{\mu}\right\rangle\right| \leq \frac{1}{2}\left|A_{1}^{1 / 2} \bar{\mu}\right|_{L^{2}}^{2} \\
+Q_{1}\left(\left|\phi_{1}\right|_{H^{1}},\left|\phi_{2}\right| H_{1^{1}}\right)\left(\left|A_{1} \phi_{1}\right|_{L^{2}}^{2}+\left|A_{1} \phi_{2}\right|_{L^{2}}^{2}\right)\|\psi\|^{2},  \tag{3.17}\\
\left|b_{1}\left(v_{2}, \psi, A_{1} \psi\right)\right| \leq \frac{\zeta}{8}\left|A_{1} \psi\right|_{L^{2}}^{2}+c\left|v_{2}\right|_{L^{2}}^{2}\left\|v_{2}\right\|^{2}\|\psi\|^{2},  \tag{3.18}\\
\zeta\left|\left\langle g_{1}\left(t, v_{1}, \phi_{1}\right)-g_{1}\left(t, v_{2}, \phi_{2}\right), w\right\rangle\right| \leq \frac{\zeta}{8}\left|A_{1} \psi\right|_{L^{2}}^{2}+c \zeta|\nabla \bar{\mu}|_{L^{2}}^{2},  \tag{3.19}\\
\leq\left. L_{1}\|w\|| |(w, \psi)\right|_{\mathcal{H}} \\
\leq \frac{1}{8}\|w\|^{2}+c L_{1}^{2}|(w, \psi)|_{\mathcal{H}}^{2},  \tag{3.20}\\
\left\|\sigma\left(s, v_{1}, \phi_{1}\right)-\sigma\left(s, v_{2}, \phi_{2}\right)\right\|_{L^{2}\left(Z, \mu, V_{1}\right)}^{2} \leq l_{2}^{2}|(w, \psi)|_{\mathcal{H}}^{2} .
\end{gather*}
$$

Let

$$
\mathcal{Y}_{2}(t)=|w(t)|_{L^{2}}^{2}+\|\psi(t)\|^{2},
$$

and

$$
\begin{gather*}
K_{1}(t)=c\left(\left\|v_{1}\right\|^{2}+\left\|\phi_{1}\right\|^{2}\left|A_{1} \phi_{1}\right|_{L^{2}}^{2}+\left\|\phi_{2}\right\|^{2}\left|A_{1} \phi_{2}\right|_{L^{2}}^{2}+\left|v_{2}\right|_{L^{2}}^{2}\left\|v_{2}\right\|^{2}\right)  \tag{3.21}\\
+Q_{1}\left(\left|\phi_{1}\right|_{H^{1}},\left|\phi_{2}\right|_{H^{1}}\right)\left(\left|A_{1} \phi_{1}\right|_{L^{2}}^{2}+\left|A_{1} \phi_{2}\right|_{L^{2}}^{2}\right)+c L_{1}^{2}+l_{2}^{2}, \\
\delta(t)=\exp \left(-\int_{0}^{t} K_{1}(s) d s\right) . \tag{3.22}
\end{gather*}
$$

Applying Itô's formula to the process $\delta(t) \mathcal{Y}_{2}(t)$ and using (3.11)-(3.20), we derive that

$$
\begin{align*}
& \mathbb{E} \delta(t) \mathcal{Y}_{2}(t)+\mathbb{E} \int_{0}^{t} \delta(s)\left(\|w\|^{2}+(1-c \zeta)\|\bar{\mu}\|^{2}+\zeta\left|A_{1} \psi\right|_{L^{2}}^{2}\right) d s \\
& \quad \leq \mathbb{E} \mathcal{Y}_{2}(0)+\mathbb{E} \int_{0}^{t} \delta(s) \mathcal{Y}_{2}(s) d s \tag{3.23}
\end{align*}
$$

Note that the expectation of the stochastic integral in (3.11) vanishes. Therefore we obtain

$$
\mathbb{E} \delta(t) \mathcal{Y}_{2}(t) \leq \mathbb{E} \mathcal{Y}_{2}(0)+\mathbb{E} \int_{0}^{t} \delta(s) \mathcal{Y}_{2}(s) d s, 0 \leq t \leq T
$$

It follows from the Gronwall lemma that there exists a constant $C>0$ such that

$$
\mathbb{E} \delta(t) \mathcal{Y}_{2}(t) \leq C \mathbb{E} \mathcal{Y}_{2}(0)
$$

for any $t \in[0, T]$, which proves the first part of the proposition. Since $\delta(t)$ is bounded and positive $\mathbb{P}$-a.s., we conclude that the second part of the proposition follows from the last estimate. Note that in (3.23), we choose $\zeta>0$ and small enough such that $1-c \zeta>0$.
3.2. Existence of solution. In this part, we prove the existence of solution. The method relies on Galerkin approximation.

Proposition 3.3. We assume that Condition 1 above is satisfied. Moreover, we suppose that $g_{1}(\cdot, 0,0) \in L^{4}\left(\Omega, L^{2}\left(0, T ; L^{2}\left(0, T ; V_{1}^{*}\right)\right),\left(v_{0}, \phi_{0}\right) \in L^{2}\left(\Omega, \mathfrak{F}_{0}, \mathbb{P} ; \mathcal{H}\right)\right.$ satisfies $\mathbb{E}\left[\mathcal{E}\left(v_{0}, \phi_{0}\right)\right]^{2}<\infty$. Then, there exists a unique solution

$$
(v, \phi) \in L^{4}\left(\Omega, \mathbb{D}\left(0, T ; H_{1}\right) \times \mathcal{C}\left(0, T ; V_{2}\right)\right) \cap L^{4}\left(\Omega, L^{2}(0, T ; \mathcal{U})\right)
$$

Furthermore, the following estimate holds:

$$
\begin{equation*}
\mathbb{E} \sup _{t \in[0, T]} \mathcal{E}(v(t), \phi(t))+\mathbb{E} \int_{0}^{T}\|(v(s), \phi(s))\|_{\mathcal{U}}^{2} d s \leq C \tag{3.24}
\end{equation*}
$$

provided that $\mathbb{E} \mathcal{E}\left(v_{0}, \phi_{0}\right)<\infty$, and

$$
\begin{equation*}
\mathbb{E} \sup _{t \in[0, T]}[\mathcal{E}(v(t), \phi(t))]^{p}+\mathbb{E}\left(\int_{0}^{T}\|(v(s), \phi(s))\|_{\mathcal{U}}^{2} d s\right)^{p} \leq C \tag{3.25}
\end{equation*}
$$

for any positive integer $p \geq 2$, provided that $\mathbb{E}\left[\mathcal{E}\left(v_{0}, \phi_{0}\right)\right]^{p}<\infty$.
Proof. Let $\left\{\left(w_{i}, \psi_{i}\right), i=1,2,3, \cdots\right\} \subset \mathcal{U}$ be an orthonormal basis of $\mathcal{H}$, where $\left\{w_{i}, i=1,2 \cdots\right\}, \quad\left\{\psi_{i}, i=1,2 \cdots\right\}$ are eigenvectors of $A_{0}$ and $A_{1}$ respectively. We set $\mathcal{U}_{m}=\mathcal{H}_{m}=\operatorname{span}\left\{\left(w_{1}, \psi_{1}\right), \cdots\left(w_{m}, \psi_{m}\right)\right\}$. We look for $\left(v_{m}, \phi_{m}\right) \in \mathcal{H}_{m}$ solution to

$$
\begin{align*}
& d v_{m}(t)=-\Pi_{m}^{1}\left[A_{0} v_{m}+B_{0}\left(v_{m}, v_{m}\right)-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right)\right] d t \\
& \quad+\Pi_{m}^{1} g_{1}\left(s, v_{m}, \phi_{m}\right) d t+\int_{Z} \Pi_{m}^{1} \sigma\left(t, v_{m}(t-), \phi_{m}(t-), z\right) \tilde{\eta}(d t, d z)  \tag{3.26}\\
& d \phi_{m}(t)=-\Pi_{m}^{2}\left[A_{1} \mu_{m}+B_{1}\left(v_{m}, \phi_{m}\right)\right] d t \\
& \mu_{m}=A_{1} \phi_{m}+f\left(\phi_{m}\right)
\end{align*}
$$

where $\Pi_{m} \equiv\left(\Pi_{m}^{1}, \Pi_{m}^{2}\right)$ is the orthogonal projection of $\mathcal{H}$ onto $\mathcal{H}_{m}$.
As in the proof of Theorem 1.2.1 of [3], we can obtain the existence and uniqueness of a solution $\left(v_{m}, \phi_{m}\right) \in L^{2}\left(\Omega \times[0, T] ; \mathcal{U}_{m}\right)$ of (3.26) on an interval $\left[0, T_{m}\right]$.

For each $n \geq 1$, we consider the $\mathfrak{F}_{t}$-stopping time $\tau_{n}$ defined by:

$$
\begin{equation*}
\tau_{n}=T \wedge \inf _{t \in[0, T]}\left\{\mathcal{E}\left(v_{m}, \phi_{m}\right)(t)+\int_{0}^{t}\left(\left\|v_{m}(s)\right\|^{2}+\left\|\bar{\mu}_{m}(s)\right\|^{2}\right) d s \geq n^{2}\right\} \tag{3.27}
\end{equation*}
$$

where hereafter $a \wedge b=\min (a, b)$.
For fixed m , the sequence $\left\{\tau_{n} ; n \geq 1\right\}$ is increasing to $T$. Throughout we fix $r \in[0, T]$ and $0 \leq t \leq r \wedge \tau_{n}$. Now using Itô's formula, we derive that as in the proof of (3.1) that

$$
\begin{align*}
\mathcal{E}\left(v_{m}, \phi_{m}\right)(t) & +\int_{0}^{t}\left(2\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s \\
& =\mathcal{E}\left(v_{0}, \phi_{0}\right)+2 \int_{0}^{t}\left\langle g_{1}\left(s, v_{m}(s), \phi_{m}(s)\right), v_{m}(s)\right\rangle d s  \tag{3.28}\\
& +2 \int_{0}^{t} \int_{Z}\left(v_{m}(s-), \sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right) \tilde{\eta}(d z, d s) \\
& +\int_{0}^{t} \int_{Z} \Upsilon(s, z) \eta(d z, d s)
\end{align*}
$$

where

$$
\begin{aligned}
\Upsilon(s, z) & =\left|v_{m}(s-)+\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right|_{L^{2}}^{2} \\
& -\left|v_{m}(s-)\right|_{L^{2}}^{2}-\left(v_{m}(s-), \sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right) .
\end{aligned}
$$

From the fact that $|x|^{2}-|y|^{2}+|x-y|^{2}=2\langle x-y, x\rangle$, it follows that

$$
\begin{align*}
\mathcal{E}\left(v_{m}(t), \phi_{m}(t)\right) & +\int_{0}^{t}\left(2\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s=\mathcal{E}\left(v_{0}, \phi_{0}\right) \\
& +2 \int_{0}^{t}\left\langle g_{1}\left(s, v_{m}(s), \phi_{m}(s)\right), v_{m}(s)\right\rangle d s  \tag{3.29}\\
& +2 \int_{0}^{t} \int_{Z}\left(v_{m}(s-), \sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right) \tilde{\eta}(d z, d s) \\
& +\int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right|_{L^{2}}^{2} \eta(d z, d s),
\end{align*}
$$

We define the following stochastic processes

$$
\begin{aligned}
X(t) & =\sup _{s \in[0, t]} \mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right) \\
Y(t) & =\int_{0}^{t}\left(2\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s \\
I(t) & =2\left|\int_{0}^{t} \int_{Z}\left(v_{m}(s-), \sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right) \tilde{\eta}(d z, d s)\right| \\
& +\int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right|_{L^{2}}^{2} \eta(d z, d s) \equiv \sup _{s \in[0, t]}\left|I_{1}(s)\right|+I_{2}(t)
\end{aligned}
$$

where

$$
\begin{align*}
& I_{1}(t)=2 \int_{0}^{t} \int_{Z}\left(v_{m}(s-), \sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right) \tilde{\eta}(d z, d s), \\
& I_{2}(t)=\sup _{s \in[0, t]} \int_{0}^{s} \int_{Z}\left|\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right|_{L^{2}}^{2} \eta(d z, d s) . \tag{3.30}
\end{align*}
$$

Since $I_{1}(t)$ is a local martingale we can apply Burkholder-Davis-Gundy's inequality to derive that

$$
\begin{align*}
& \mathbb{E} \sup _{s \in\left[0, r \wedge \tau_{n}\right]}\left|I_{1}(s)\right| \leq \\
& \quad C \mathbb{E}\left(\int_{0}^{r \wedge \tau_{n}} \int_{Z}\left(v_{m}(s-), \sigma\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)\right)^{2} \nu(d z) d s\right)^{1 / 2} . \tag{3.31}
\end{align*}
$$

Thanks to Hölder's and Young's inequalities we have

$$
\begin{align*}
\mathbb{E} \sup _{s \in[0, t]}\left|I_{1}(s)\right| & \leq C\left[\epsilon \mathbb{E} \sup _{s \in[0, t]}\left|v_{m}(s)\right|_{L^{2}}^{2}\right]^{1 / 2} \times \\
& {\left[\epsilon^{-1} \mathbb{E} \int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right|_{L^{2}}^{2} \nu(d z) d s\right]^{1 / 2} } \\
& \leq C \epsilon \mathbb{E} \sup _{s \in[0, t]}\left|v_{m}(s)\right|_{L^{2}}^{2}  \tag{3.32}\\
& +C \epsilon^{-1} \mathbb{E} \int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right|_{L^{2}}^{2} \nu(d z) d s \\
& \leq C \epsilon \mathbb{E} \sup _{s \in[0, t]} \mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right) \\
& +C \epsilon^{-1} \mathbb{E} \int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right|_{L^{2}}^{2} \nu(d z) d s .
\end{align*}
$$

Using (2.27), we derive that

$$
\begin{equation*}
\mathbb{E} \sup _{s \in[0, t]}\left|I_{1}(s)\right| \leq C \epsilon X(t)+C \epsilon^{-1} l_{0} t+C \epsilon^{-1} \int_{0}^{t} \mathbb{E} X(s) d s . \tag{3.33}
\end{equation*}
$$

Next, we will deal with the second term of $I(t)$. Taking into account that the process

$$
\int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)\right|^{2} \eta(d z, d s)
$$

has only positive jumps, we derive from (2.27) that

$$
\begin{align*}
& \mathbb{E} I_{2}(t) \leq \mathbb{E} \int_{0}^{t} \int_{Z}\left|\sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right|_{L^{2}}^{2} \nu(d z) d s \\
& \leq l_{0} t+l_{1} \int_{0}^{t} \mathbb{E} \mid\left(v_{m}(s),\left.\phi_{m}(s)\right|_{\mathcal{H}} ^{2} d s\right.  \tag{3.34}\\
& \leq l_{0} t+l_{1} \int_{0}^{t} \mathbb{E} X(s) d s
\end{align*}
$$

We also have

$$
\begin{align*}
\left|2\left\langle g_{1}\left(s, v_{m}, \phi_{m}\right), v_{m}\right\rangle\right| & \leq 2 L_{1}\left|\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}\left\|v_{m}\right\|+2\left\|g_{1}(s, 0,0)\right\|_{V_{1}^{*}}\left\|v_{m}\right\| \\
& \leq \frac{1}{8}\left\|v_{m}\right\|^{2}+c L_{1}^{2}\left|\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}^{2}+c\left\|g_{1}(s, 0,0)\right\|_{V_{1}^{*}}^{2} \tag{3.35}
\end{align*}
$$

It follows from (3.28)-(3.35) that

$$
\begin{align*}
\mathbb{E} \mathcal{E}\left(v_{m}, \phi_{m}\right)(t) & +\mathbb{E} \int_{0}^{t}\left(\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s \leq \mathbb{E} \mathcal{E}\left(v_{0}, \phi_{0}\right) \\
& +c \mathbb{E} \int_{0}^{t} \mathcal{E}\left(v_{m}, \phi_{m}\right)(s) d s+c \mathbb{E} \int_{0}^{t}\left\|g_{1}(s, 0,0)\right\|_{V_{1}^{*}}^{2} d s \tag{3.36}
\end{align*}
$$

Therefore from Lemma 2.5, we derive that there exist a positive constant $C$ such that

$$
\begin{equation*}
\mathbb{E} \mathcal{E}\left(v_{m}, \phi_{m}\right)(t)+\mathbb{E} \int_{0}^{t}\left(\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s \leq C \tag{3.37}
\end{equation*}
$$

for any $m \in \aleph$ and $t \in\left[0, r \wedge \tau_{n}\right], r \in[0, T]$.
We have just shown that $\forall t \in[0, T]$, we have

$$
\begin{equation*}
\mathbb{E} \sup _{s \in\left[0, t \wedge \tau_{n}\right]} \mathcal{E}\left(v_{m}, \phi_{m}\right)(t)+\mathbb{E} \int_{0}^{t}\left(\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s \leq C \tag{3.38}
\end{equation*}
$$

from which we can infer that

$$
\mathbb{P}\left(\tau_{n}<t\right) \leq C n^{-2}, \forall t \in[0, T], \forall n>0
$$

Hence, $\lim _{n \rightarrow+\infty} \mathbb{P}\left(\tau_{n}<t\right)=0$, for all $t \in[0, T]$. That is, $\tau_{n} \rightarrow+\infty$ in probability. Therefore, there exists a subsequence $\tau_{n_{k}}$ such that $\tau_{n_{k}} \rightarrow+\infty$, a.s. Since the sequence $\left(\tau_{n}\right)_{n}$ is increasing, we infer that $\tau_{n_{k}} \nearrow+\infty$ a.s.. Now we use Fatou's lemma and pass to the limit in (3.38) and derive that

$$
\begin{equation*}
\mathbb{E} \sup _{s \in[0, t]} \mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)+\mathbb{E} \int_{0}^{t}\left(\left\|v_{m}(s)\right\|^{2}+\left\|\bar{\mu}_{m}(s)\right\|^{2}\right) d s \leq C \tag{3.39}
\end{equation*}
$$

and (3.24) is proved.
To prove (3.25), we proceed as follows. First, without loss of generality, we may assume that for any given $m$, the process $\left(v_{m}, \phi_{m}\right)$ is uniformly bounded in $[0, T]$. Otherwise, we can introduce a sequence of stopping times as before.

By raising both sides of (3.28) to the power of $p \geq 2$, we derive that

$$
\begin{align*}
& \mathbb{E} \sup _{s \in[0, t]}\left[\mathcal{E}\left(v_{m}(t), \phi_{m}(t)\right)\right]^{p}+\mathbb{E}\left[\int_{0}^{t}\left(2\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s\right]^{p} \\
& \quad \leq \mathbb{E}\left[\mathcal{E}\left(v_{0}, \phi_{0}\right)\right]^{p}+c \mathbb{E} \sup _{s \in[0, t]}\left|I_{3}(s)\right|^{p}+c \mathbb{E} \sup _{s \in[0, t]}\left|I_{4}(s)\right|^{p}  \tag{3.40}\\
& \quad+c \mathbb{E}\left[\int_{0}^{t}\left\|g_{1}(s, 0,0)\right\|_{V_{1}^{*}}^{2} d s\right]^{p}
\end{align*}
$$

where

$$
\begin{align*}
& I_{3}(t)= \\
& \int_{0}^{t} \int_{Z}\left\{\left|v_{m}(s-)+\Pi_{m}^{1} \sigma\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)\right|_{L^{2}}^{2}-\left|v_{m}(s-)\right|_{L^{2}}^{2}\right\} \tilde{\eta}(d z, d s) \tag{3.41}
\end{align*}
$$

$$
\begin{align*}
I_{4}(t) & = \\
& \left.\left.\int_{0}^{t} \int_{Z}\left\{\mid v_{m}(s)+\Pi_{m}^{1} \sigma\left(s, v_{m}, \phi_{m}\right)(s), z\right)\right|_{L^{2}} ^{2}-\left|v_{m}(s)\right|_{L^{2}}^{2}\right\} \nu(d z) d s \\
& -2 \int_{0}^{t} \int_{Z}\left(v_{m}(s), \Pi_{m}^{1} \sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right) \nu(d z) d s  \tag{3.42}\\
& \left.\leq c \int_{0}^{t} \int_{Z} \mid \Pi_{m}^{1} \sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right)\left.\right|_{L^{2}} ^{2} \nu(d z) d s \\
& \leq c \int_{0}^{t}\left(1+\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2}\right) d s
\end{align*}
$$

As in $[10,13]$, we note that

$$
\begin{align*}
& \int_{Z}\left\{\left|v_{m}(s-)+\Pi_{m}^{1} \sigma\left(s,\left(v_{m}, \phi_{m}\right)(s), z\right)\right|_{L^{2}}^{2}-\left|v_{m}(s-)\right|_{L^{2}}^{2}\right\}^{2} \nu(d z) \\
& \quad \leq\left|v_{m}(s-)\right|_{L^{2}}^{2} \int_{Z}\left|\sigma\left(s,\left(v_{m}, \phi_{m}\right)(s), z\right)\right|_{L^{2}}^{2} \nu(d z) \\
& \quad+c \int_{Z}\left|\sigma\left(s,\left(v_{m}, \phi_{m}\right)(s), z\right)\right|_{L^{2}}^{4} \nu(d z)  \tag{3.43}\\
& \quad \leq c_{0}+c_{1}\left|v_{m}(s)\right|_{L^{2}}^{2}+c_{2}\left|v_{m}(s)\right|_{L^{2}}^{4} \\
& \quad \leq k_{1}+k_{4}\left|v_{m}(s)\right|_{L^{2}}^{4}
\end{align*}
$$

It follows that

$$
\begin{align*}
& \left(\int_{0}^{t} \int_{Z}\left\{\left|v_{m}(s-)+\Pi_{m}^{1} \sigma\left(s, v_{m}, \phi_{m}, z\right)\right|_{L^{2}}^{2}-\left|v_{m}(s-)\right|_{L^{2}}^{2}\right\}^{2} \nu(d z) d s\right)^{p / 2}  \tag{3.44}\\
& \quad \leq c\left(k_{1} T\right)^{p / 2}+c\left(k_{2}\right)^{p / 2}\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{4} d s\right)^{p / 2}
\end{align*}
$$

We derive that

$$
\begin{align*}
\mathbb{E} \sup _{s \in[0, t]} & \left|I_{3}(s)\right|^{p} \leq c_{p}\left(k_{1} T\right)^{p / 2} \\
& +c_{p}\left(k_{2}\right)^{p / 2} \mathbb{E}\left[\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{4} d s\right)^{p / 2}\right]  \tag{3.45}\\
& \leq c+\frac{1}{2} \mathbb{E}\left(\sup _{s \in[0, t]}\left|v_{m}(s)\right|_{L^{2}}^{2}\right)^{p}+c \mathbb{E}\left(\int_{0}^{t}\left|v_{m}(s)\right|_{L^{2}}^{2} d s\right)^{p} .
\end{align*}
$$

From Hölder's inequality, we have

$$
\begin{align*}
\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2} d s & \leq\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2 p} d s\right)^{1 / p}\left(\int_{0}^{t} 1 d s\right)^{\frac{p-1}{p}}  \tag{3.46}\\
& \leq T^{\frac{p-1}{p}}\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2 p} d s\right)^{1 / p}
\end{align*}
$$

which gives

$$
\begin{equation*}
\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2} d s\right)^{p} \leq c T^{P-1} \int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2 p} d s \tag{3.47}
\end{equation*}
$$

From (3.45), (3.46), we get

$$
\begin{align*}
\mathbb{E} & \sup _{s \in[0, t]}\left|I_{3}(s)\right|^{p} \leq \frac{1}{2} \mathbb{E}\left(\sup _{s \in[0, t]}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2}\right)^{p} \\
& +c_{p, T} \int_{0}^{t} \mathbb{E}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2 p} d s  \tag{3.48}\\
& \leq \frac{1}{2} \mathbb{E} \sup _{s \in[0, t]}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p}+c_{p, T} \int_{0}^{t} \mathbb{E}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p} d s .
\end{align*}
$$

From (2.27) and (3.42), we also have

$$
\begin{align*}
\mathbb{E}\left|I_{4}(t)\right|^{p} & \leq c \mathbb{E}\left(\int_{0}^{t}\left(1+\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2}\right) d s\right)^{p}  \tag{3.49}\\
& \left.\leq c_{p}+c_{p} \mathbb{E}\left(\int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2}\right) d s\right)^{p}
\end{align*}
$$

It follows that

$$
\begin{align*}
\mathbb{E} \sup _{s \in[0, t]}\left|I_{2}(s)\right|^{p} & \leq c_{p, T}+c_{p, T} \int_{0}^{t}\left|\left(v_{m}, \phi_{m}\right)(s)\right|_{\mathcal{H}}^{2 p} d s  \tag{3.50}\\
& \leq c_{p, T}+c_{p, T} \int_{0}^{t} \mathbb{E}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p} d s .
\end{align*}
$$

It follows from (3.40)-(3.50)

$$
\begin{align*}
\mathbb{E} & \sup _{s \in[0, t]}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p} \leq c_{p, T} \\
& +c_{p, T} \int_{0}^{t} \mathbb{E}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p} d s+c\left(\int_{0}^{t}\left\|g_{1}(t, 0,0)\right\|_{V_{1}^{*}}^{2}\right)^{p}, \tag{3.51}
\end{align*}
$$

and Gronwall's lemma and (3.40) give

$$
\begin{equation*}
\mathbb{E} \sup _{s \in[0, t]}\left[\mathcal{E}\left(v_{m}(s), \phi_{m}(s)\right)\right]^{p}+\mathbb{E}\left[\int_{0}^{t}\left(2\left\|v_{m}(s)\right\|^{2}+\left\|\mu_{m}(s)\right\|^{2}\right) d s\right]^{p} \leq C, \tag{3.52}
\end{equation*}
$$

and (3.25) follows.
Proposition 3.4. We can extract from $\left(v_{m}, \phi_{m}\right)$ a subsequence still labeled the same and there exists a stochastic process $(v, \phi)$ such that

$$
\begin{aligned}
& \left(v_{m}, \phi_{m}\right) \rightharpoonup(v, \phi) \text { in } L^{4}\left(\Omega, L^{\infty}([0, T] ; \mathcal{H})\right), \\
& \left(v_{m}, \phi_{m}\right) \rightharpoonup(v, \phi) \text { in } L^{2}\left(\Omega, L^{2}([0, T] ; \mathcal{U})\right), \\
& B_{0}\left(v_{m}, v_{m}\right) \rightharpoonup \beta_{0}^{b}, R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right) \rightharpoonup r_{0}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right), \\
& g_{1}\left(t, v_{m}, \phi_{m}\right) \rightharpoonup g_{1}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right), \\
& B_{1}\left(v_{m}, \phi_{m}\right) \rightharpoonup \beta_{1}^{b}, f\left(\phi_{m}\right) \rightharpoonup f^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{2}^{*}\right) \text {, } \\
& \sigma\left(t, v_{m}, \phi_{m}, \cdot\right) \rightharpoonup \sigma^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; L^{2}\left(Z, \nu ; H_{1}\right)\right) \text {. }
\end{aligned}
$$

We note that

$$
\begin{align*}
& \left|B_{0}\left(v_{m}, v_{m}\right)\right|_{V_{1}^{*}} \leq c\left|v_{m}\right|_{L^{2}}\left\|v_{m}\right\|, \\
& \left|R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right)\right|_{V_{1}^{*}} \leq c\left\|\phi_{m}\right\|\left|A_{1} \phi_{m}\right|_{L^{2}}^{1 / 2}\left|\phi_{m}\right|_{H^{3}}^{1 / 2}  \tag{3.54}\\
& \left|B_{1}\left(v_{m}, \phi_{m}\right)\right|_{V_{2}^{*}} \leq c\left|v_{m}\right|_{L^{2}}\left\|\phi_{m}\right\|^{1 / 2}\left|A_{1} \phi_{m}\right|_{L^{2}}^{1 / 2}
\end{align*}
$$

From (2.11) and $(3.26)_{3}$, we also have

$$
\begin{align*}
& \left|f\left(\phi_{m}\right)\right|_{L^{2}} \leq c\left(1+\left\|\phi_{m}\right\|^{k+1}\right) \leq c\left(1+\mathcal{E}\left(v_{m}, \phi_{m}\right)^{\frac{k+1}{2}}\right), \\
& \left|A_{1} \phi_{m}\right|_{L^{2}} \leq c\left|\mu_{m}\right|_{L^{2}}+c\left|f\left(\phi_{m}\right)\right|_{L^{2}}  \tag{3.55}\\
& \quad \leq c\left|\mu_{m}\right|_{L^{2}}+c\left(1+\mathcal{E}\left(v_{m}, \phi_{m}\right)^{\frac{k+1}{2}}\right), \\
& \quad\left|\phi_{m}\right|_{H^{3}}^{2} \quad \leq c\left\|\mu_{m}\right\|^{2}+c\left|f^{\prime}\left(\phi_{m}\right) \nabla \phi_{m}\right|_{L^{2}}^{2}  \tag{3.56}\\
& \quad \leq c\left\|\mu_{m}\right\|^{2}+c\left\|\phi_{m}\right\|^{2 k+1}\left|A_{1} \phi_{m}\right|_{L^{2}} .
\end{align*}
$$

It follows from (3.54)-(3.55) that

$$
\begin{align*}
& \mathbb{E} \sup _{[0, T]}\left|\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}^{2} \leq C, \quad \mathbb{E} \int_{0}^{T} \|\left(v_{m}(s), \phi_{m}(s) \|_{\mathcal{U}}^{2} d s \leq C,\right. \\
& \mathbb{E} \int_{0}^{T}\left|\phi_{m}(s)\right|_{H^{3}}^{2} d s \leq C, \quad \mathbb{E} \sup _{t \in[0, T]}\left|f\left(\phi_{m}\right)\right|_{L^{2}} \leq C, \\
& \mathbb{E} \int_{0}^{T}\left[\left|B_{0}\left(v_{m}, v_{m}\right)\right|_{V_{1}^{*}}^{2}+\left|R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right)\right|_{V_{1}^{*}}^{2}+\left|B_{1}\left(v_{m}, \phi_{m}\right)\right|_{V_{2}^{*}}^{2}\right] d s \leq C,  \tag{3.57}\\
& \mathbb{E} \sup _{[0, T]}\left|\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}^{4} \leq C, \quad \mathbb{E}\left[\int_{0}^{T} \|\left(v_{m}(s-), \phi_{m}(s-) \|_{\mathcal{U}}^{2} d s\right]^{2} \leq C\right. \\
& \mathbb{E} \int_{0}^{T}\left\|\sigma\left(s, v_{m}(s), \phi_{m}(s), z\right)\right\|_{L^{2}\left(Z, \nu, H_{1}\right)}^{2} d s \leq l_{0} T  \tag{3.58}\\
& \quad+l_{1} \mathbb{E} \int_{0}^{T}\left|\left(v_{m}(s-), \phi_{m}(s-)\right)\right|_{\mathcal{H}}^{2} d s \leq C \\
& \mathbb{E} \int_{0}^{T}\left|f\left(\phi_{m}\right)\right|_{L^{2}}^{2} d s \leq c \mathbb{E} \int_{0}^{T}\left(1+\mathcal{E}\left(v_{m}, \phi_{m}\right)^{k+1}\right) d s \leq C . \tag{3.59}
\end{align*}
$$

From (3.57), we can find a subsequence still denoted $\left\{\left(v_{m}, \phi_{m}\right)\right\}$ such that

$$
\begin{align*}
& \left(v_{m}, \phi_{m}\right) \rightharpoonup(v, \phi) \text { in } L^{4}\left(\Omega, L^{\infty}([0, T] ; \mathcal{H})\right), \\
& \left(v_{m}, \phi_{m}\right) \rightharpoonup(v, \phi) \text { in } L^{2}(\Omega \times[0, T] ; \mathcal{U}), \\
& B_{0}\left(v_{m}, v_{m}\right) \rightharpoonup \beta_{0}^{b}, R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right) \rightharpoonup r_{0}^{b}, \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right), \\
& g_{1}\left(t, v_{m}, \phi_{m}\right) \rightharpoonup g_{1}^{b}, f\left(\phi_{m}\right) \rightharpoonup f^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right),  \tag{3.60}\\
& B_{1}\left(v_{m}, \phi_{m}\right) \rightharpoonup \beta_{1}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{2}^{*}\right), \\
& \sigma\left(t, v_{m}(s-), \phi_{m}(s-), \cdot\right) \rightharpoonup \sigma^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; L^{2}\left(Z, \nu ; H_{1}\right)\right) .
\end{align*}
$$

As in $[11,29]$, we can check that $v$ is an $H_{1}$-valued càdlàg and $\mathbb{F}$-progressively measurable process, and $\phi$ is an $V_{2}$-valued continuous and $\mathbb{F}$-progressively measurable process. Moreover $(v, \phi)$ satisfies for all $0 \leq t \leq T$

$$
\begin{align*}
v(t) & +\int_{0}^{t} A_{0} v d s+\int_{0}^{t} \beta_{0}^{b}(s) d s=v_{0}+\int_{0}^{t}\left(r_{0}^{b}(s)+g_{1}^{b}(s)\right) d s \\
& +\int_{0}^{t} \int_{Z} \sigma^{b}(s, z) \tilde{\eta}(d z, d s)  \tag{3.61}\\
\phi(t) & +\int_{0}^{t} A_{1} \mu^{b} d s+\int_{0}^{t} \beta_{1}^{b}(s) d s=\phi_{0}, \mu^{b}=A_{1} \phi+f^{b},
\end{align*}
$$

$\mathbb{P}$-a.s. as a equality in $\mathcal{U}^{*}$.

Proposition 3.5. We have the following identities

$$
\begin{align*}
& \beta_{0}^{b}=B_{0}(v, v), r_{0}^{b}=R_{0}\left(A_{1} \phi, \phi\right)  \tag{3.62}\\
& \beta_{1}^{b}=B_{1}(v, \phi), f^{b}=f(\phi), \sigma(t, v, \phi, z)=\sigma^{b}
\end{align*}
$$

Proposition 3.6. For any $n \geq 1$ we have that as $m \longrightarrow+\infty$,

$$
\begin{equation*}
1_{\left[0, \tau_{n}\right]}\left(\left(v_{m}, \phi_{m}\right)-(v, \phi)\right) \longrightarrow(0,0) \text { in } L^{2}(\Omega \times[0, T] ; \mathcal{U}) \tag{3.63}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{E}\left|\left(v_{m}, \phi_{m}\left(\tau_{n}\right)\right)-(v, \phi)\left(\tau_{n}\right)\right|_{\mathcal{H}} \longrightarrow 0 \text { as } n \longrightarrow+\infty . \tag{3.64}
\end{equation*}
$$

Proof. Let

$$
\left(\tilde{v}_{m}, \tilde{\phi}_{m}, \tilde{\mu}_{m}\right)=\Pi_{m}(v, \phi, \mu)
$$

where $\Pi_{m} \equiv\left(\Pi_{m}^{1}, \Pi_{m}^{2}\right)$ is the orthogonal projection of $\mathcal{H}$ onto $\mathcal{H}_{m}$. It follows that

$$
\begin{align*}
& \left|\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)\right|_{\mathcal{H}} \leq|(v, \phi)|_{\mathcal{H}} \\
& \left.\left\|\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)\right\|_{\mathcal{U}} \leq c \| v, \phi\right) \|_{\mathcal{U}}, \\
& \left(\tilde{v}_{m}, \tilde{\phi}_{m}\right) \rightarrow(v, \phi) \text { in } \mathcal{U} \text { for almost every }(\omega, t) \in \Omega \times[0, T],  \tag{3.65}\\
& \left(\tilde{v}_{m}, \tilde{\phi}_{m}\right) \rightarrow(v, \phi) \text { in } L^{2}(\Omega \times[0, T] ; \mathcal{U}), \\
& \mathbb{E}\left(\left|\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)\left(\tau_{n}\right)-(v, \phi)\left(\tau_{n}\right)\right|_{\mathcal{H}}\right) \rightarrow 0 \text { as } n \rightarrow+\infty .
\end{align*}
$$

From (3.26) and (3.61), we derive that for $1 \leq k \leq m$, we have

$$
\begin{align*}
& \left\langle\tilde{v}_{m}(t)-v_{m}(t), w_{k}\right\rangle+\int_{0}^{t}\left\langle A_{0}\left(\tilde{v}_{m}-v_{m}\right), w_{k}\right\rangle d s \\
& \quad+\int_{0}^{t}\left\langle\beta_{0}^{b}-B_{0}\left(v_{m}, v_{m}\right), w_{k}\right\rangle d s=\int_{0}^{t}\left\langle r_{0}^{b}-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right), w_{k}\right\rangle d s \\
& \quad+\int_{0}^{t}\left\langle g_{1}^{b}-g_{1}\left(s, v_{m}, \phi_{m}\right), w_{k}\right\rangle d s \\
& \quad+\int_{0}^{t} \int_{Z}\left[\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)-\sigma^{b}(s, z), w_{k}\right] \tilde{\eta}(d z, d s)  \tag{3.66}\\
& \left\langle\tilde{\phi}_{m}(t)-\phi_{m}(t), A_{1} \psi_{k}\right\rangle+\int_{0}^{t}\left\langle A_{1}\left(\tilde{\mu}_{m}-\mu_{m}\right), A_{1} \psi_{k}\right\rangle d s \\
& \quad+\int_{0}^{t}\left\langle\beta_{1}^{b}-B_{1}\left(v_{m}, \phi_{m}\right), A_{1} \psi_{k}\right\rangle d s=0 \\
& \left\langle\tilde{\mu}_{m}-\mu_{m}, A_{1} \psi_{k}\right\rangle=\left\langle A_{1}\left(\tilde{\phi}_{m}-\phi_{m}\right), A_{1} \psi_{k}\right\rangle+\left\langle f^{b}-f\left(\phi_{m}\right), A_{1} \psi_{k}\right\rangle .
\end{align*}
$$

Note that

$$
\begin{aligned}
& \beta_{0}^{b}-B_{0}\left(v_{m}, v_{m}\right)=\beta_{0}^{b}-B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right)+B_{0}\left(\tilde{v}_{m}-v_{m}, \tilde{v}_{m}\right) \\
& \quad \quad+B_{0}\left(v_{m}, \tilde{v}_{m}-v_{m}\right), \\
& r_{0}^{b}-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right)=r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right)+R_{0}\left(A_{1}\left(\tilde{\phi}_{m}-\phi_{m}\right), \tilde{\phi}_{m}\right) \\
& \quad \quad+R_{0}\left(A_{1} \phi_{m}, \phi_{m}-\phi_{m}\right), \\
& \beta_{1}^{b}-B_{1}\left(v_{m}, \phi_{m}\right)=\beta_{1}^{b}-B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)+B_{1}\left(\tilde{v}_{m}-v_{m}, \tilde{\phi}_{m}\right) \\
& \quad+B_{1}\left(v_{m}, \tilde{\phi}_{m}-\phi_{m}\right), \\
& f^{b}-f\left(\phi_{m}\right)=f^{b}-f\left(\tilde{\phi}_{m}\right)+f\left(\tilde{\phi}_{m}\right)-f\left(\phi_{m}\right) .
\end{aligned}
$$

Let us set $\theta_{m}=\tilde{v}_{m}-v_{m}, \rho_{m}=\phi_{m}-\phi_{m}, \zeta_{m}=\tilde{\mu}_{m}-\mu_{m}$. From the Itô's formula, we have

$$
\begin{align*}
& \left|\theta_{m}(t)\right|_{L^{2}}^{2}+2 \int_{0}^{t}\left(\left\|\theta_{m}\right\|^{2}+\left\langle\beta_{0}^{b}-B_{0}\left(v_{m}, v_{m}\right), \theta_{m}\right\rangle\right) d s \\
& \quad=2 \int_{0}^{t}\left\langle r_{0}^{b}-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right), \theta_{m}\right\rangle d s+2 \int_{0}^{t}\left\langle g_{1}^{b}-g_{1}\left(s, v_{m}, \phi_{m}\right), \theta_{m}\right\rangle d s  \tag{3.67}\\
& \quad+2 \int_{0}^{t} \int_{Z}\left[\sigma\left(v_{m}(s-), \phi_{m}(s-), z\right)-\Pi_{m}^{1} \sigma^{b}(s, z)\right] \tilde{\eta}(d z, d s) \\
& \quad+\int_{0}^{t} \int_{Z} \Upsilon(s, z) \eta(d z, d s),
\end{align*}
$$

where

$$
\begin{aligned}
& \Upsilon(s, z)=\left|v_{m}(s-)+\sigma_{m}\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)-\Pi_{m}^{1} \sigma^{b}(s, z)\right|_{L^{2}}^{2} \\
& \quad-\left|v_{m}(s-)\right|_{L^{2}}^{2}-2\left\langle\sigma_{m}\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)-\Pi_{m}^{1} \sigma^{b}(s, z), v_{m}(s-)\right\rangle \\
& \quad=\left|\sigma\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)-\Pi_{m}^{1} \sigma^{b}(s, z)\right|_{L^{2}}^{2} .
\end{aligned}
$$

Replacing $\psi_{k}$ in $(3.66)_{3}$ and $(3.66)_{2}$ respectively by $\bar{\zeta}_{m}-\xi \rho_{m}$ and $\rho_{m}$ gives

$$
\begin{align*}
& \frac{d}{d t}\left\|\rho_{m}\right\|^{2}+2\left\|\bar{\zeta}_{m}\right\|^{2}+2 \xi\left|A_{1} \rho_{m}\right|_{L^{2}}^{2}+2\left\langle\beta_{1}^{b}-B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle \\
& \quad+2\left\langle B_{1}\left(\theta_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle+2\left\langle B_{1}\left(v_{m}, \rho_{m}\right), A_{1} \rho_{m}\right\rangle+\xi\left\langle\zeta_{m}, A_{1} \rho_{m}\right\rangle  \tag{3.68}\\
& \quad+\xi\left\langle f^{b}-f\left(\phi_{m}\right), A_{1} \rho_{m}\right\rangle-\left\langle f^{b}-f\left(\phi_{m}\right), A_{1} \zeta_{m}\right\rangle=0 .
\end{align*}
$$

Note that

$$
\begin{align*}
& \left\langle\beta_{0}^{b}-B_{0}\left(v_{m}, v_{m}\right), \theta_{m}\right\rangle=\left\langle\beta_{0}^{b}-B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle+\left\langle B_{0}\left(\theta_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle \\
& \leq\left\langle\beta_{0}^{b}-B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle+\frac{1}{4}\left\|\theta_{m}\right\|^{2}+c\left\|\tilde{v}_{m}\right\|^{2}\left|\theta_{m}\right|_{L^{2}}^{2},  \tag{3.69}\\
& \left\langle\beta_{1}^{b}-B_{1}\left(v_{m}, \phi_{m}\right), A_{1} \rho_{m}\right\rangle=\left\langle\beta_{1}^{b}-B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle \\
& +\left\langle B_{1}\left(\theta_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle+\left\langle B_{1}\left(v_{m}, \rho_{m}\right), A_{1} \rho_{m}\right\rangle \\
& \leq\left\langle\beta_{1}^{b}-B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle+\frac{1}{4}\left(\left\|\theta_{m}\right\|^{2}+\frac{\xi}{4}\left|A_{1} \rho_{m}\right|_{L^{2}}^{2}\right)  \tag{3.70}\\
& +c\left\|\phi_{m}\right\|^{2}\left|A_{1} \phi_{m}\right|_{L^{2}}^{2}\left|\theta_{m}\right|_{L^{2}}^{2}+c\left|\tilde{v}_{m}\right|_{L^{2}}^{2}\left\|\tilde{v}_{m}\right\|^{2}\left\|\rho_{m}\right\|^{2}, \\
& \left\langle r_{0}^{b}-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right), \theta_{m}\right\rangle=\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle \\
& +\left\langle R_{0}\left(A_{1} \rho_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle+\left\langle R_{0}\left(A_{1} \phi_{m}, \rho_{m}\right), \theta_{m}\right\rangle \\
& \leq\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle+\frac{1}{4}\left(\left\|\theta_{m}\right\|^{2}+\frac{\xi}{4}\left|A_{1} \rho_{m}\right|_{L^{2}}^{2}\right)  \tag{3.71}\\
& +c\left|A_{1} \phi_{m}\right|_{L^{2}}^{2}\left(\left|\theta_{m}\right|_{L^{2}}^{2}+\left\|\rho_{m}\right\|^{2}\right)+c\left\|\tilde{\phi}_{m}\right\|^{2}\left|A_{1} \tilde{\phi}_{m}\right|_{L^{2}}^{2}\left|\theta_{m}\right|_{L^{2}}^{2}, \\
& \left\langle f^{b}-f\left(\phi_{m}\right), A_{1} \zeta_{m}\right\rangle=\left\langle f^{b}-f\left(\tilde{\phi}_{m}\right), A_{1} \zeta_{m}\right\rangle+\left\langle f\left(\tilde{\phi}_{m}\right)-f\left(\phi_{m}\right), A_{1} \zeta_{m}\right\rangle \\
& \leq\left\langle f^{b}-f\left(\tilde{\phi}_{m}\right), A_{1} \zeta_{m}\right\rangle+\frac{1}{2}\left\|\zeta_{m}\right\|^{2}  \tag{3.72}\\
& +Q_{1}\left(\left\|\tilde{\phi}_{m}\right\|,\left\|\phi_{m}\right\|\right)\left(\left|A_{1} \phi_{m}\right|_{L^{2}}^{2}+\left|A_{1} \tilde{\phi}_{m}\right|_{L^{2}}^{2}\right)\left\|\rho_{m}\right\|^{2}, \\
& \xi\left\langle f^{b}-f\left(\phi_{m}\right), A_{1} \rho_{m}\right\rangle=\xi\left\langle f^{b}-f\left(\tilde{\phi}_{m}\right)+f\left(\tilde{\phi}_{m}\right)-f\left(\phi_{m}\right), A_{1} \rho_{m}\right\rangle \\
& \leq \xi\left\langle f^{b}-f\left(\tilde{\phi}_{m}\right), A_{1} \zeta_{m}\right\rangle+\frac{\xi}{8}\left|A_{1} \rho_{m}\right|_{L^{2}}^{2}  \tag{3.73}\\
& +Q_{1}\left(\left\|\tilde{\phi}_{m}\right\|,\left\|\phi_{m}\right\|\right)\left\|\rho_{m}\right\|^{2} \text {, } \\
& \left.\left|\xi\left\langle\zeta_{m}, A_{1} \rho_{m}\right\rangle \leq \frac{\xi}{2}\left\|\zeta_{m}\right\|^{2}+\frac{\xi}{4}\right| A_{1} \rho_{m}\right|_{L^{2}} ^{2},  \tag{3.74}\\
& \Upsilon(s, z)=\left|\sigma_{m}\left(s, v_{m}(s-), \phi_{m}(s-), z\right)-\Pi_{m}^{1} \sigma^{b}(s, z)\right|_{L^{2}}^{2} \\
& =\mid \Pi_{m}^{1}\left[\sigma\left(s, v_{m}(s-), \phi_{m}(s-), z\right)-\left.\sigma(s, v(s-), \phi(s-), z)\right|_{L^{2}} ^{2}\right.  \tag{3.75}\\
& -\mid \Pi_{m}^{1}\left[\sigma(s, v(s-), \phi(s-), z)-\left.\sigma^{b}(s, z)\right|_{L^{2}} ^{2}+S_{1}(s, z),\right.
\end{align*}
$$

where

$$
\begin{aligned}
& S_{1}(s, z)= \\
& 2\left\langle\Pi_{m}^{1}\left[\sigma\left(s,\left(v_{m}, \phi_{m}\right)(s-), z\right)-\sigma^{\mathrm{b}}(s, z)\right], \Pi_{m}^{1}\left[\sigma(s,(v, \phi)(s-), z)-\sigma^{\mathrm{b}}(s, z)\right]\right\rangle .
\end{aligned}
$$

From (2.27) and (3.75), we derive that

$$
\begin{align*}
\Upsilon(s, z) & \leq l_{2}\left|\left(v_{m}(s-), \phi_{m}(s-)\right)\right|_{\mathcal{H}}^{2}+l_{2}\left|\left(v_{m}, \phi_{m}\right)(s-)-(v, \phi)(s-)\right|_{\mathcal{H}}^{2}  \tag{3.76}\\
& -\mid \Pi_{m}^{1}\left[\sigma(s, v(s-), \phi(s-), z)-\left.\sigma^{b}(s, z)\right|_{L^{2}} ^{2}+S_{1}(s, z)\right.
\end{align*}
$$

We also have

$$
\begin{align*}
\left\langle g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right)-g_{1}\left(s, v_{m}, \phi_{m}\right), \theta_{m}\right\rangle & \leq L_{1}\left|\left(\theta_{m}, \rho_{m}\right)\right|_{\mathcal{H}}\left\|\theta_{m}\right\|  \tag{3.77}\\
& \leq \frac{1}{4}\left\|\theta_{m}\right\|^{2}+c L_{1}^{2}\left|\left(\theta_{m}, \rho_{m}\right)\right|_{\mathcal{H}}^{2}
\end{align*}
$$

Let

$$
\begin{aligned}
Z(t) & =\left|\theta_{m}(t)\right|_{L^{2}}^{2}+\left\|\rho_{m}(t)\right\|^{2}=\left|\left(\tilde{v}_{m}-v_{m}\right)(t)\right|_{L^{2}}^{2}+\left\|\left(\tilde{\phi}_{m}-\phi_{m}\right)(t)\right\|^{2} \\
Y_{1}(t) & \left.=c\left\|\tilde{v}_{m}\right\|^{2}+c\left\|\phi_{m}\right\|^{2}\left|A_{1} \phi_{m}\right|_{L^{2}}^{2}+c\left|\tilde{v}_{m}\right|_{L^{2}}^{2}\left\|\tilde{v}_{m}\right\|^{2}+\left.c\left|A_{1} \phi_{m}\right|_{L^{2}}^{2} \tilde{x}^{2}\right|^{2}\right) \\
& +c\left\|\tilde{\phi}_{m}\right\|^{2}\left|A_{1} \tilde{\phi}_{m}\right|_{L^{2}}^{2}+Q_{1}\left(\left\|\tilde{\phi}_{m}\right\|,\left\|\phi_{m}\right\|\right)\left(1+\left|A_{1} \phi_{m}\right|_{L^{2}}^{2}+\left|A_{1} \tilde{\phi}_{m}\right|_{L^{2}}\right) \\
& +c L_{2}^{2}+c L_{1}^{2} \\
K_{2}(t) & =\left\|\theta_{m}\right\|^{2}+(1-c \xi)\left\|\zeta_{m}\right\|^{2}+c \xi\left|A_{1} \rho_{m}\right|_{L^{2}}^{2},
\end{aligned}
$$

where $\xi$ is small enough such that $1-c \xi>0$.
Let

$$
\delta(t)=\exp \left(-\int_{0}^{t} Y_{1}(s) d s\right)
$$

Using (3.67)-(3.77), it follows from Ito's formula that

$$
\begin{align*}
\mathbb{E} \delta(t) Z(t) & +\mathbb{E} \int_{0}^{t} \delta(s) K_{2}(s) d s \\
& +\mathbb{E} \int_{0}^{t} \delta(s) \mid \Pi_{m}^{1}\left[\sigma(s,(v, \phi)(s), z)-\left.\sigma^{b}(s, z)\right|_{L^{2}} ^{2} d s\right. \\
& \leq \mathbb{E} \int_{0}^{t} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s \\
& +\mathbb{E} \int_{0}^{t} \delta(s)\left\langle-\beta_{1}^{b}+B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle d s  \tag{3.78}\\
& +\mathbb{E} \int_{0}^{t} \delta(s)\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s \\
& +\mathbb{E} \int_{0}^{t} \delta(s)\left\langle g_{1}^{b}-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}\right), \theta_{m}\right\rangle d s \\
& +\mathbb{E} \int_{0}^{t} \int_{Z} \delta(s) S_{1}(s, z) \eta(d z, d s) .
\end{align*}
$$

For each $n \geq 1$, we consider the $\mathfrak{F}_{t}-$ stopping time $\tau_{n}$ defined by:

$$
\tau_{n}=\min \left(T, \inf \left\{t \in[0, T] ;|(v, \phi)|_{\mathcal{H}}^{2}+\int_{0}^{t}\|(v, \phi)\|_{\mathcal{U}}^{2} d s \geq n^{2}\right\}\right)
$$

We derive from (3.78) that

$$
\begin{align*}
\mathbb{E}\left(\tau_{n}\right) \delta(s) & Z\left(\tau_{n}\right)+c \mathbb{E} \int_{0}^{\tau_{n}} \delta(s) K_{2}(s) d s \\
& +\mathbb{E} \int_{0}^{\tau_{n}} \delta(s) \mid \Pi_{m}^{1}\left[\sigma(s,(v, \phi)(s), z)-\left.\sigma^{b}(s, z)\right|_{L^{2}} ^{2} d s\right. \\
& \leq 2 \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s \\
& +2 \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}^{b}-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s  \tag{3.79}\\
& +2 \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s \\
& +2 \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{1}^{b}+B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle d s \\
& +\mathbb{E} \int_{0}^{Q_{n}} \int_{Z} \delta(s) S_{1}(s, z) \eta(d z, d s) .
\end{align*}
$$

Claim 1. The right side of (3.79) goes to 0 as $m$ goes to $+\infty$.
(i). Since $\Pi_{m}^{1} \circ \Pi_{m}^{1}=\Pi_{m}^{1}$ and $\left\|\Pi_{m}^{1}\right\| \leq 1$, it follows that

$$
1_{\left[0, \tau_{n}\right]} \delta(s) \Pi_{m}^{1}\left[\sigma(s, v(s), \phi(s), z)-\sigma^{b}(s, z)\right]
$$

is bounded in $L^{2}\left(\Omega \times[0, T] ; L^{2}\left(Z, \nu ; H_{1}\right)\right)$. Therefore, from $(3.60)_{6}$ we see that

$$
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \int_{Z} \delta(s) S_{1}(s, z) \eta(d z, d s)=0
$$

(ii). Let us now prove that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.80}
\end{equation*}
$$

We recall that

$$
\begin{align*}
& \left(v_{m}, \phi_{m}\right) \rightharpoonup(v, \phi),\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right) \rightarrow(v, \phi),\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)-\left(v_{m}, \phi_{m}\right) \rightharpoonup(0,0)  \tag{3.81}\\
& \text { in } L^{2}(\Omega \times[0, T] ; \mathcal{U}) .
\end{align*}
$$

We also have

$$
\begin{align*}
\lim _{m \rightarrow \infty} & \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s \\
& =\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}(v, v), \theta_{m}\right\rangle d s  \tag{3.82}\\
& +\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-B_{0}(v, v)+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s
\end{align*}
$$

From (3.81) and the fact that $\left.1\right|_{\left[0, \tau_{n}\right]} \delta(t)\left(-\beta_{0}^{b}+B_{0}(v, v)\right) \in L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right)$, it follows that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}(v, v), \theta_{m}\right\rangle d s=0 \tag{3.83}
\end{equation*}
$$

We also note that

$$
\begin{aligned}
& \left\|-B_{0}(v, v)+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right)\right\|_{V_{1}^{*}} \\
& \quad \leq c\left|\tilde{v}_{m}-v\right|_{L^{2}}^{1 / 2}\left\|\tilde{v}_{m}-v\right\|^{1 / 2}\left(\left\|\tilde{v}_{m}\right\|^{1 / 2}\left|\tilde{v}_{m}\right|_{L^{2}}^{1 / 2}+\|v\|^{1 / 2}|v|_{L^{2}}^{1 / 2}\right)
\end{aligned}
$$

which implies that

$$
\begin{aligned}
& \left\|\left.1\right|_{\left[0, \tau_{n}\right]}\left(-B_{0}(v, v)+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right)\right)\right\|_{V_{1}^{*}} \rightarrow 0 \text { as } m \rightarrow \infty, d t \times d \mathbb{P}-\text { a.e., } \\
& \left\|\left.1\right|_{\left[0, \tau_{n}\right]}\left(-B_{0}(v, v)+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right)\right)\right\|_{V_{1}^{*}} \leq c n\|v\| \in L^{2}(\Omega \times[0, T] ; \Re) .
\end{aligned}
$$

It follows that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-B_{0}(v, v)+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.84}
\end{equation*}
$$

We conclude from (3.83) and (3.84) that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{0}^{b}+B_{0}\left(\tilde{v}_{m}, \tilde{v}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.85}
\end{equation*}
$$

which proves (3.80).
(iii). Next we will prove that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.86}
\end{equation*}
$$

From (3.81) and the fact that

$$
\left.1\right|_{\left[0, \tau_{n}\right]} \delta(t)\left(r_{0}^{b}-R_{0}\left(A_{1} \phi, \phi\right)\right) \in L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right),
$$

we also have

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle r_{0}^{b}-R_{0}\left(A_{1} \phi, \phi\right), \theta_{m}\right\rangle d s=0 \tag{3.87}
\end{equation*}
$$

We also note that

$$
\begin{aligned}
\left\|R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right)-R_{0}\left(A_{1} \phi, \phi\right)\right\|_{V_{1}^{*}} & \leq c\left|A_{1}\left(\tilde{\phi}_{m}-\phi\right)\right|_{L^{2}}\|\phi\|^{1 / 2}\left|A_{1} \phi\right|_{L^{2}}^{1 / 2} \\
& +c\left\|\tilde{\phi}_{m}-\phi\right\|^{1 / 2}\left|A_{1}\left(\tilde{\phi}_{m}-\phi\right)\right|_{L^{2}}^{1 / 2}\left|A_{1} \tilde{\phi}_{m}\right|_{L^{2}}
\end{aligned}
$$

which implies that

$$
\begin{aligned}
& \left\|\left.1\right|_{\left[0, \tau_{n}\right]}\left(R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right)-R_{0}\left(A_{1} \phi, \phi\right)\right)\right\|_{V_{1}^{*}} \rightarrow 0 \text { as } m \rightarrow \infty, d t \times d \mathbb{P}-\text { a.e. }, \\
& \left\|\left.1\right|_{\left[0, \tau_{n}\right]}\left(R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right)-R_{0}\left(A_{1} \phi, \phi\right)\right)\right\|_{V_{1}^{*}} \leq c n\|v\| \in L^{2}(\Omega \times[0, T] ; \Re) .
\end{aligned}
$$

It follows that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right)-R_{0}\left(A_{1} \phi, \phi\right), \theta_{m}\right\rangle d s=0 \tag{3.88}
\end{equation*}
$$

We conclude from (3.87) and (3.88) that

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle r_{0}^{b}-R_{0}\left(A_{1} \tilde{\phi}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.89}
\end{equation*}
$$

which proves (3.86).
(iv). Let us now prove that

$$
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{1}^{b}+B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle d s=0
$$

Following similar steps as in (3.85) and (3.89), can check that

$$
\begin{aligned}
\lim _{m \rightarrow \infty} & \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{1}^{b}+B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle d s \\
& =\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-\beta_{1}^{b}+B_{1}(v, \phi), A_{1} \rho_{m}\right\rangle d s \\
& +\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle-B_{1}(v, \phi)+B_{1}\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right), A_{1} \rho_{m}\right\rangle d s=0
\end{aligned}
$$

$(\mathbf{v})$. Let us also prove that

$$
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}^{b}(s)-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s=0
$$

From (3.81) and the fact that

$$
\begin{aligned}
& 1 \mid\left[0, \tau_{n}\right] \delta(t)\left(g_{1}^{b}(t)-g_{1}(t, v, \phi)\right) \in L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right) \text { and } \\
& 1 \mid\left[0, \tau_{n}\right] \delta(t)\left(g_{1}(t, v, \phi)-g_{1}\left(t, \tilde{v}_{m}, \tilde{\phi}_{m}\right)\right) \rightarrow 0 \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right) \\
& \text { as } m \rightarrow \infty,
\end{aligned}
$$

we derive that

$$
\begin{align*}
& \lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}^{b}(s)-g_{1}(s, v, \phi), \theta_{m}\right\rangle d s=0 \\
& \lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}(s, v, \phi)-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s=0 \tag{3.90}
\end{align*}
$$

Therefore, we derive that

$$
\begin{aligned}
\lim _{m \rightarrow \infty} & \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}^{b}(s)-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s \\
& =\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}^{b}(s)-g_{1}(s, v, \phi), \theta_{m}\right\rangle d s \\
& +\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \delta(s)\left\langle g_{1}(s, v, \phi)-g_{1}\left(s, \tilde{v}_{m}, \tilde{\phi}_{m}\right), \theta_{m}\right\rangle d s=0
\end{aligned}
$$

The proof of the convergence of the other terms is similar.
Finally we conclude that the right side of (3.79) goes to 0 as $m$ goes to $+\infty$.
Now using the fact that $1 \mid\left[0, \tau_{n}\right] \delta(t) \leq 1$, we derive from (3.79) that

$$
\begin{align*}
& \lim _{m \rightarrow \infty} \mathbb{E}\left|\left(\theta_{m}, \psi_{m}\right)\left(\tau_{n}\right)\right|_{\mathcal{H}}^{2} \\
& \quad=\lim _{m \rightarrow \infty} \mathbb{E}\left(\left|\left(\tilde{v}_{m}, \tilde{\phi}_{m}\right)\left(\tau_{n}\right)-\left(v_{m}, \phi_{m}\right)\left(\tau_{n}\right)\right|_{\mathcal{H}}^{2}\right)=0 \\
& \lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} K_{2}(s) d s \\
& \quad=\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}}\left(\left\|\theta_{m}\right\|^{2}+(1-c \xi)\left\|\zeta_{m}\right\|^{2}+\xi\left|A_{1} \rho_{m}\right|_{L^{2}}^{2}\right) d s  \tag{3.91}\\
& \quad=\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}}\left(\left\|\tilde{v}_{m}-v_{m}\right\|^{2}+(1-c \xi)\left\|\tilde{\mu}_{m}-\mu_{m}\right\|^{2}\right) d s \\
& \quad+\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}} \xi\left|A_{1}\left(\tilde{\phi}_{m}-\phi_{m}\right)\right|_{L^{2}}^{2} d s=0 .
\end{align*}
$$

We now give the proof of Proposition 3.4.
Proof of Proposition 3.4. Our goal is to prove that the following hold true.

$$
\begin{align*}
& \sigma(s, v, \phi, z)=\sigma^{b}(s, z) \text { in } L^{2}\left(\Omega \times[0, T] ; L^{2}\left(Z, \nu, H_{1}\right)\right), \\
& B_{0}(v, v)=\beta_{0}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right) \\
& R_{0}\left(A_{1} \phi, \phi\right)=r_{0}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right),  \tag{3.92}\\
& B_{1}(v, \phi)=\beta_{1}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{2}^{*}\right), \\
& g_{1}(t, v, \phi)=g_{1}^{b}(t) \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right) .
\end{align*}
$$

It is clear that $(3.92)_{1}$, follows from $(3.53)_{6}$.
To prove $(3.92)_{2}$, we proceed as follows. We note that from $(3.91)_{2}$ and (3.65), we also have

$$
\left.\left.\left(v_{m}, \phi_{m}\right)\right|_{\left[0, \tau_{n}\right]} \rightarrow(v, \phi)\right|_{\left[0, \tau_{n}\right]} \text { in } L^{2}(\Omega \times[0, T] ; \mathcal{U})
$$

Therefore, for any $w \in L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)$, we have

$$
\begin{aligned}
& \mathbb{E} \int_{0}^{\tau_{n}}\left\langle B_{0}(v, v)-B_{0}\left(v_{m}, v_{m}\right), w\right\rangle d s \leq \\
& c\|w\|_{L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)} \times \mathbb{E} \int_{0}^{\tau_{n}}\left\|v_{m}-v\right\|^{1 / 2}\left|v_{m}-v\right|_{L^{2}}^{1 / 2}\left(\|v\|+\left\|v_{m}\right\|\right) d s
\end{aligned}
$$

which gives

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}}\left\langle B_{0}(v, v)-B_{0}\left(v_{m}, v_{m}\right), w\right\rangle d s=0 \tag{3.93}
\end{equation*}
$$

From $(3.60)_{3}$ and (3.93), we derive that

$$
\mathbb{E} \int_{0}^{\tau_{n}}\left\langle B_{0}(v, v)-\beta_{0}^{b}, w\right\rangle=0, \forall w \in L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)
$$

Since $\tau_{n} \uparrow T$ and $L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)$ is dense in $L^{2}\left(\Omega \times[0, T] ; V_{1}\right)$, we conclude that

$$
B_{0}(v, v)=\beta_{0}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right) .
$$

This proves $(3.92)_{2}$.
To prove $(3.92)_{3}$, we note that

$$
\begin{aligned}
& \mathbb{E} \int_{0}^{\tau_{n}}\left\langle R_{0}\left(A_{1} \phi, \phi\right)-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right), w\right\rangle d s \leq \\
& c\|w\|_{L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)} \mathbb{E} \int_{0}^{\tau_{n}}\left|A_{1}\left(\phi_{m}-\phi\right)\right|_{L^{2}}^{1 / 2}\|\phi\|^{1 / 2}\left|A_{1} \phi\right|_{L^{2}}^{1 / 2} d s+ \\
& c\|w\|_{L^{\infty}\left(\Omega \times[0, T] ; V_{1}\right)} \mathbb{E} \int_{0}^{\tau_{n}}\left\|\phi_{m}-\phi\right\|^{1 / 2}\left|A_{1}\left(\phi_{m}-\phi\right)\right|_{L^{2}}^{1 / 2}\left|A_{1} \phi_{m}\right|_{L^{2}} d s,
\end{aligned}
$$

which gives

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{\tau_{n}}\left\langle R_{0}\left(A_{1} \phi, \phi\right)-R_{0}\left(A_{1} \phi_{m}, \phi_{m}\right), w\right\rangle d s=0 \tag{3.94}
\end{equation*}
$$

From $(3.60)_{4}$ and (3.94), we derive that

$$
\mathbb{E} \int_{0}^{\tau_{n}}\left\langle R_{0}\left(A_{1} \phi, \phi\right)-r_{0}^{b}, w\right\rangle=0, \forall w \in L^{2}\left(\Omega \times[0, T] ; V_{1}\right)
$$

which gives

$$
R_{0}\left(A_{1} \phi, \phi\right)=r_{0}^{b} \text { in } L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right)
$$

and $(3.92)_{3}$ is proved.
Similarly, we can prove that
$B_{1}(v, \phi)=\beta_{1}^{b}$ in $L^{2}\left(\Omega \times[0, T] ; V_{2}^{*}\right), g_{1}(t, v, \phi)=g_{1}^{b}(t)$ in $L^{2}\left(\Omega \times[0, T] ; V_{1}^{*}\right)$.
3.3. Convergence of the whole sequence of the Galerkin approximation. In this part, we prove the convergence of the whole sequence of the Galerkin approximation to the solution $(v, \phi)$ of (2.24).

Theorem 3.7. The whole sequence of solutions to the Galerkin approximation $\left\{\left(v_{m}, \phi_{m}\right) ; m \in \aleph\right\}$ defined by (3.26) satisfies

$$
\begin{align*}
& \lim _{m \rightarrow \infty} \mathbb{E}\left|\left(v_{m}, \phi_{m}\right)(T-)-(v, \phi)(T-)\right|_{\mathcal{H}}^{2}=0 \\
& \lim _{m \rightarrow \infty} \mathbb{E}\left\|\left(v_{m}, \phi_{m}\right)-(v, \phi)\right\|_{\mathcal{U}}^{2}=0 \tag{3.95}
\end{align*}
$$

Proof. For the proof, we first recall from $[9,16]$ the following lemma.
Lemma 3.8. Let $\left\{Q_{m} ; m \geq 1\right\} \subset L^{2}(\Omega \times[0, T] ; \Re)$ be a sequence of continuous real processes, and let $\left\{\tau_{n} ; n \geq 1\right\}$ be a sequence of $\mathfrak{F}_{t}$-stopping times such that $\tau_{n} \uparrow T$; $\sup _{m>1} \mathbb{E}\left|Q_{m}(T)\right|^{2}<\infty$, and $\lim _{m \rightarrow \infty} \mathbb{E}\left|Q_{m}\left(\tau_{n}\right)\right|=0$, for $n \geq 1$. Then $\lim _{m \rightarrow \infty} \mathbb{E}\left|Q_{m}(T)\right|=$ 0.

Applying Lemma 3.8 to $Q_{m}(t)=\left|(v, \phi)-\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}^{2}$ and $\delta_{n}=\tau_{n}$ and using (3.25), (3.91) $)_{1}$ and the uniqueness of $(v, \phi)$, we conclude that the whole sequence given by (3.26) satisfies

$$
\lim _{m \rightarrow \infty} \mathbb{E}\left|(v, \phi)-\left(v_{m}, \phi_{m}\right)\right|_{\mathcal{H}}^{2}=0, \forall t \in[0, T]
$$

Similarly, applying Lemma 3.8 to $Q_{m}(t)=\int_{0}^{t}\left\|(v, \phi)(s)-\left(v_{m}, \phi_{m}\right)(s)\right\|_{\mathcal{U}}^{2} d s$ and using (3.25), $(3.91)_{2}$, we conclude that the whole sequence $\left(v_{m}, \phi_{m}\right)$ converges to $(v, \phi)$ strongly in $L^{2}(\Omega \times[0, T] ; \mathcal{U})$, i.e.,

$$
\lim _{m \rightarrow \infty} \mathbb{E} \int_{0}^{t}\left\|(v, \phi)(s)-\left(v_{m}, \phi_{m}\right)(s)\right\|_{\mathcal{U}}^{2} d s=0, \quad \forall t \in[0, T] .
$$

Acknowledgments. The research of the first author is supported by the Fulbright Scholar Program Advanced Research and Florida International University.

## References

1. Abels, H.: On a diffuse interface model for a two-phase flow of compressible viscous fluids. Indiana Univ. Math. J., 57:659-698, 2008.
2. Abels, H.: On a diffuse interface model for two-phase flows of viscous, incompressible fluids with matched densities. Arch. Ration. Mech. Anal., 194:463-506, 2009.
3. Albeverio, S., Brzeźniak, Z. and Wu, J. L.: Existence of global solutions and invariant measures for stochastic differential equations driven by poisson type noise with non-Lipschitz coefficients. J. Math. Anal. Appl., 371(1):309-322, 2010.
4. Bensoussan, A. and Temam, R.: Equations stochastiques de type Navier-Stokes. Journal of Functional Analysis, 13:195-222, 1973.
5. Blesgen, T.: A generalization of the Navier-Stokes equation to two-phase flow. Pysica $D$ (Applied Physics), 32:1119-1123, 1999.
6. Boyer, F.: Mathematical study of multi-phase flow under shear through order parameter formulation. Asymptotic Anal., 20:175-212, 1999.
7. Boyer, F.: Nonhomogeneous Cahn-Hilliard fluids. Ann. Inst. H. Poincaré Anal. Non Linéaire, 18:225-259, 2001.
8. Boyer, F.: A theoretical and numerical model for the study of incompressible mixture flows. Computer and Fluids, 31:41-68, 2002.
9. Breckner (Lisei), H.: Approximation and optimal control of the stochastic Navier-Stokes equations. Dissertation, Martin-Luther University, Halle-Wittenberg, 1999.
10. Brzeźniak, Z., Hausenblas, E. and Razafimandimby, P. A.: Some results on the penalised nematic liquid crystals driven by multiplicative noise. arXiv:1310.8641, 2016.
11. Brzeźniak, Z., Hausenblas, E. and Zhu, J.: 2D stochastic Navier-Stokes equations driven by jump noise. Nonlinear Anal., 79:122-139, 2013.
12. Brzeźniak, Z., Liu, W. and Zhu, J.: Strong solutions for SPDE with locally monotone coefficients driven by Lévy noise. Nonlinear Anal. Real World Appl., 17:283-310, 2014.
13. Brzeźniak, Z., Manna, U. and Panda, A. A.: Existence of weak martingale solution of nematic liquid crystals driven by pure jump noise. arXiv:1706.05056, 2017.
14. Caginalp, G.: An analysis of a phase field model of a free boundary. Arch. Rational Mech. Anal., 92(3):205-245, 1986.
15. Cao, C. and Gal, G.: Global solutions for the 2D NS-CH model for a two-phase flow of viscous, incompressible fluids with mixed partial viscosity and mobility. Nonlinearity, 25(11):3211-3234, 2012.
16. Caraballo, T., Real, J. and Taniguchi, T.: On the existence and uniqueness of solutions to stochastic three-dimensional Lagrangian averaged Navier-Stokes equations. Proc. R. Soc. Lond. Ser. A Math. Phys. Eng. Sci., 462(2066):459-479, 2006.
17. Chueshov, I. and Millet, A.: Stochastic 2D hydrodynamical type systems: well posedness and large deviations. Appl. Math. Optim., 61:379-420, 2010.
18. Deugoué, G. and Tachim Medjo, T.: The exponential behavior of a stochastic globally modified Cahn-Hilliard-Navier-Stokes model with multiplicative noise. J. Math. Anal. Appl., 460(1):140-163, 2018.
19. Deugoué, G. and Tachim Medjo, T.: On the convergence for the 3D globally modified Cahn-Hilliard-Navier-Stokes equations.tachimt@fiu.edu J. Differential Equations, 265(2):545592, 2018.
20. Dong, Z., Li, W. V. and Zhai, J.: Stationary weak solutions for stochastic 3D Navier-Stokes equations with Lévy noise. Stochastic and Dynamics, 12(1):1150006, 2012.
21. Dong, Z. and Xie, Y.: Global solutions of stochastic 2D Navier-Stokes equations with Lévy noise. Science in China Series A: Mathematics, 52(7):1497-1524, 2009.
22. Dong, Z. and Zhai, J.: Martingale solutions and Markov selection of stochastic 3D NavierStokes equations with jump. Journal of Differential Equations, 250(6):2737-2778, 2011.
23. Faris, W. G. and Jona-Lasinio, G.: Large fluctuations for a nonlinear heat equation with noise. J. Phys. A, 15(10):3025-3055, 1982.
24. Feireisl, E., Petzeltová, E., Rocca, E. and Schimperna, G.: Analysis of a phase-field model for two-phase compressible fluids. Math. Models Methods Appl. Sci., 20(7):1129-1160, 2010.
25. Gal, C. and Grasselli, M.: Asymptotic behavior of a Cahn-Hilliard-Navier-Stokes system in 2D. Ann. Inst. H. Poincaré Anal. Non Linéaire, 27(1):401-436, 2010.
26. Gal., C. and Grasselli, M.: Longtime behavior for a model of homogeneous incompressible two-phase flows. Discrete Contin. Dyn. Syst., 28(1):1-39, 2010.
27. Gal, C. and Grasselli, M.: Trajectory attractors for binary fluid mixtures in 3D. Chin. Ann. Math. Ser. B, 31(5):655-678, 2010.
28. Gurtin, M. E., Polignone, D. and Vinals, J.: Two-phase binary fluid and immiscible fluids described by an order parameter. Math. Models Methods Appl. Sci., 6:8-15, 1996.
29. Hausenblas., E. and Razafimandimby, P. A.: On stochastic evolution equations for nonlinear bipolar fluids: well-posedness and some properties of the solution. J. Math. Anal. Appl., 441(2):763-800, 2016.
30. Hohenberg, P. C. and Halperin, B. I.: Theory of dynamical critical phenomena. Rev. Modern Phys., 49:435-479, 1977.
31. Leray, J.: Sur le mouvement d'un liquide visqueux emplissant l'espace. Acta Math., 63(1):193248, 1934.
32. Lions, J. L. and Prodi, G.: Un théorème d'existence et unicité dans les équations de NavierStokes en dimension 2. C. R. Acad. Sci. Paris, 248:3519-3521, 1959.
33. Liu, W. and Röckner, M.: SPDE in Hilbert space with locally monotone coefficients. $J$. Funct. Anal., 259(11):2902-2922, 2010.
34. Liu, W. and Röckner, M.: Local and global well-posedness of SPDE with generalized coercivity conditions. J. Differential Equations, 254(2):725-755, 2013.
35. Tachim Medjo, T.: On the existence and uniqueness of solution to a stochastic 2D Cahn-Hilliard-Navier-Stokes model. J. Differential Equations, 262:1028-1054, 2017.
36. Motyl, E.: Martingale solution to the 2D and 3D stochastic Navier-Stokes equations driven by the compensated poisson random measure. Department of Mathematics and Computer Sciences, Lodz University, Preprint 13, 2011.
37. Motyl, E.: Stochastic Navier-Stokes equations driven by Lévy noise in unbounded 3D domains. Potential Analysis, 38:863-912, 2013.
38. Motyl, E.: Stochastic hydrodynamic-type evolution equations driven by Lévy noise in 3D unbounded domains- Abstract framework and applications. Stochastic Processes and their Applications, 124:2052-2097, 2014.
39. Peszat, S. and Zabczyk, J.: Stochastic partial differential equations with Lévy noise. An evolution equation approach. Encyclopedia of Mathematics and its Applications, 113. Cambridge University Press, Cambridge, 2007.
40. Razafimandimby, P. A., Hausenblas, E. and Sango, M.: Martingale solution to equations for differential type fluids of grade two driven by random force of Lévy type. Potential Analysis, 38(4):1291-1331, 2013.
41. Temam, R.: Infinite Dynamical Systems in Mechanics and Physics, volume 68. Appl. Math. Sci., Springer-Verlag, New York, second edition, 1997.
42. Temam, R.: Navier-Stokes Equations, Theory and Numerical Analysis. AMS-Chelsea Series, AMS, Providence, 2001.
43. Temam, R. and Ziane, M.: Some mathematical problems in geophysical fluid dynamics. In S. Friedlander and D. Serre, editors, Handbook of Mathematical Fluid Dynamics, Vol. III, pages 535-658. Elsevier, 2004.
44. Xu, T. and Zhang, T.: Large deviation principles for 2D stochastic Navier-Stokes equations driven by Lévy processes. Journal of Functional Analysis, 257(5):1519-1545, 2009.
45. Zhai, J. and Zhang, T.: Large deviations for 2D stochastic Navier-Stokes equations driven by multiplicative Lévy noises. Bernoulli, 21(4):2351-239, 2015.
46. Mustafa Kemal Yilmaz, Sidika Basci \& Muruvvet Buyukboyaci., Earnings Yield, Market-to-Book Ratio and Dividend Yield as Predictors of Stock Returns: Evidence from the Istanbul Stock Exchange, The Global Journal of Finance and Economics,
47. Subarna K. Samanta \& Sascha Szyfman., Impact of an Interest on Reserves Regime on Monetary Policy Effectiveness: Evidence from New Zealand, Journal of World Economic Review
48. Yu-Tung Peng, Hue Hwa Au Yong and Sirimon Treepongkaruna, Contagion and Flight-to-Quality: Evidences from the Asia-Pacific Economic Cooperation (APEC) Region, Review of Applied Economics
49. Valerio Lintner., The Political Economy of European Economic Integration and the Economic Sovereignty of the Nations State, International Journal of Economics
50. Shuang-Ling Chong \& Ken P. Chong, Chong Cycle, Durability Issues and Accelerated Tests of Bridge Coatings, Journal of Mechanics and MEMS
51. Bassim Shebeb \& Ashraf Nakibullah, Understanding Monetary Policy of Bahrain, Global Review of Business and Economic Research
52. Anjana Jain \& KalyaniVijayan: X-ray Determination of the Thermal Expansion Coefficient of Nylon 6,6 Fibers, International Journal of Electrospun Nanofibers and Applications
53. Eiber Albrecht: On the Simulation of Human Hearing, International Journal of Computational Vision and Biomechanics
54. Song-Hao ZHU and Yun-Cai LIU, Automatic Video Partition for High-Level Search, International Journal of Computer Science, Information Technology and Management
55. Elina Mikelsone: IDEA MANAGEMENT AND ORGANISATIONAL EFFECTIVENESS: A RESEARCH GAP, Journal of Business Management
56. Magdi S. Mahmoud, Delay-Dependent Robust Stability and Control of Uncertain Discrete Singular Systems with State-Delay, International Journal of Contemporary Mathematics
57. Houman B. Rokni, Reconstruction of Variational Iteration Method for Boundary Value Problems in Structural Engineering and Fluid Mechanics, International Journal of Nonlinear Dynamics in Engineering and Sciences
58. A. H. Khater, Chebyshev Solution of Integral Equations with Singular Kernel, International Journal of Computational Mathematics and Numerical Simultation
59. Jian-Wei LIU, Research on Traffic Data-Collecting System Based on MC9S12D64 Microcontroller, Journal of Information Technology and Engineering
60. William P. Brown, Paul P. Eggermont, Vincent N. LaRiccia and Roland R. Roth, Partioned spline estimators for growth curve estimation in wildlife studies, International Journal of Statistics and Management System
61. Nikolaos Sarrianidis, Evangelos Drimbetas \& George Konteos., Impact of International Volatility, the Euro, and Derivatives on a Small Stock Market, International Review of Applied Economic Research
62. Sidney Askew and Agatha E. Jeffers, The Proposed Expansion in Auditors Responsibilities Author, International Journal Data Modelling and Knowledge Management
63. Fakhri, M., Farzaneh, M.R., Eslamian, S. and Khordadi, M.J. (2012) 'Uncertainty assessment of downscaled rainfall: impact of climate change on the probability of flood', Journal of Flood Engineering
64. Ratna Vadra., State Level Fiscal Reforms Agenda in India, Journal of Asian Business Management
65. Takashi Obara., Fluctuation-Dissipation Theorems on Foreign Exchange Data, International Journal of Financial Economics and Econometrics
G. Deugoué: DEpartment of Mathematics, Florida International University, DM413B, University Park, Miami, Florida 33199, USA AND:
Department of Mathematics and Computer Science, University of Dschang, P.O. Box 67, DSCHANG, CAMEROON

E-mail address: agdeugoue@yahoo.fr
T. Tachim Medjo: Department of Mathematics, Florida International University, DM413B, University Park, Miami, Florida 33199, USA

E-mail address: tachimt@fiu.edu


[^0]:    Received 2019-3-13; Accepted 2019-6-13; Communicated by the editors.
    2010 Mathematics Subject Classification. 35R60,35Q35,60H15,76M35,86A05.
    Key words and phrases. Stochastic; Cahn-Hilliard, Navier-Sokes; strong solutions; Lévy noise.

    * Corresponding author.

