# Asymptotic behavior of a semilinear problem in heat conduction with long time memory and non-local diffusion 

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#### Abstract

In this paper, the asymptotic behavior of a semilinear heat equation with long time memory and non-local diffusion is analyzed in the usual set-up for dynamical systems generated by differential equations with delay terms. This approach is different from the previous published literature on the long time behavior of heat equations with memory which is carried out by the Dafermos transformation. As a consequence, the obtained results provide complete information about the attracting sets for the original problem, instead of the transformed one. In particular, the proved results also generalize and complete previous literature in the local case.


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## 1. Introduction

The main objective of this paper is to analyze the asymptotic behavior of a semilinear heat equation with long time memory and non-local diffusion, which is an interesting situation with important applications in the real world.

On the one hand, the effects that memory terms (or the past history of a phenomenon) produce on the evolution of a dynamical system is obvious, since it is sensible to think that the evolution of any system depends not only on the current state but on its whole history (see, for instance, [1, 8, 12, 2, 6, 10, 15] and the references therein). On the other hand, many problems are better described by considering non-local terms, which created a great interest in the modeling of various real applications (see [3, 4, 5, 12] and the references therein).

[^0]Motivated by some physical problems from thermal memory or materials with memory, one can find a significant literature devoted to the analysis of partial differential equations with long time memory. For example, the authors introduced in [12] a semilinear partial differential equation to model the heat flow in a rigid, isotropic, homogeneous heat conductor with linear memory, which is given by

$$
\begin{cases}c_{0} \partial_{t} u-k_{0} \Delta u-\int_{-\infty}^{t} k(t-s) \Delta u(s) d s+f(u)=h, & \text { in } \Omega \times(\tau,+\infty)  \tag{1.1}\\ u(x, t)=0, & \text { on } \partial \Omega \times(\tau,+\infty) \\ u(x, \tau+t)=u_{0}(x, t), & \text { in } \Omega \times(-\infty, 0]\end{cases}
$$

where $\Omega \subset \mathbb{R}^{N}$ is a bounded domain with regular boundary, $u: \Omega \times \mathbb{R} \rightarrow \mathbb{R}$ is the temperature field, $k: \mathbb{R}^{+} \rightarrow \mathbb{R}$ is the heat flux memory kernel, $\mathbb{R}^{+}$denotes the interval $(0,+\infty), c_{0}$ and $k_{0}$ denote the specific heat and the instantaneous conductivity, respectively. To solve (1.1) successfully, the authors considered this problem as a non-delay one by making the past history of $u$ from $-\infty$ to $0^{-}$be part of the forcing term given by the causal function $g$, which is defined by

$$
g(x, t)=h(x, t)+\int_{-\infty}^{\tau} k(t-s) \Delta u_{0}(x, s) d s, \quad x \in \Omega, \quad t \geq \tau
$$

In this way, 1.1 becomes an initial value problem without delay or memory,

$$
\begin{cases}c_{0} \partial_{t} u-k_{0} \Delta u-\int_{\tau}^{t} k(t-s) \Delta u(s) d s+f(u)=g, & \text { in } \Omega \times(\tau,+\infty)  \tag{1.2}\\ u(x, t)=0, & \text { on } \partial \Omega \times(\tau,+\infty) \\ u(x, \tau)=u_{0}(x, 0) & \text { in } \Omega\end{cases}
$$

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However, this problem does not generate a dynamical system in an appropriate phase space, since the equation in $(1.2)$ depends on the past history and we are just fixing an initial value at time $\tau$.

Therefore, two alternatives are possible. The first one is based on the idea introduced by Dafermos [7], for linear viscoelasticity, in the 70's. Let us define the new variables,

$$
u^{t}(x, s)=u(x, t-s), \quad s \geq 0, \quad t \geq \tau
$$

$$
\begin{equation*}
\eta^{t}(x, s)=\int_{0}^{s} u^{t}(x, r) d r=\int_{t-s}^{t} u(x, r) d r, \quad s \geq 0, \quad t \geq \tau \tag{1.3}
\end{equation*}
$$

Besides, assuming $k(\infty)=0$, a change of variable and a formal integration by parts imply

$$
\int_{-\infty}^{t} k(t-s) \Delta u(s) d s=-\int_{0}^{\infty} k^{\prime}(s) \Delta \eta^{t}(s) d s
$$

Setting

$$
\mu(s)=-k^{\prime}(s),
$$

the original equation (1.2) turns into the following autonomous system without delay,

$$
\begin{cases}c_{0} \frac{\partial u}{\partial t}-k_{0} \Delta u-\int_{0}^{\infty} \mu(s) \Delta \eta^{t}(s) d s+f(u)=g, & \text { in } \Omega \times(\tau, \infty),  \tag{1.4}\\ \eta_{t}^{t}(s)=-\eta_{s}^{t}(s)+u(t), & \text { in } \Omega \times(\tau, \infty) \times \mathbb{R}^{+}, \\ u(x, t)=\eta^{t}(x, s)=0, & \text { on } \partial \Omega \times \mathbb{R} \times \mathbb{R}^{+}, \\ u(x, \tau)=u_{0}(0), & \text { in } \Omega, \\ \eta^{\tau}(x, s)=\eta_{0}(s), & \text { in } \Omega \times \mathbb{R}^{+},\end{cases}
$$

where, $\eta_{s}^{t}$ denotes the distributional derivative of $\eta^{t}(s)$ with respect to the internal variable $s$. It follows from the definition of $\eta^{t}(x, s)$ (see (1.3)) that

$$
\begin{equation*}
\eta_{0}(s)=\int_{\tau-s}^{\tau} u(r) d r=\int_{\tau-s}^{\tau} u_{0}(r-\tau) d r=\int_{-s}^{0} u_{0}(r) d r \tag{1.5}
\end{equation*}
$$

which is the initial integrated past history of $u$ with vanishing boundary. Consequently, any solution to (1.2) is a solution to (1.4) for the corresponding initial values $\left(u_{0}(0), \eta_{0}\right)$ given by (1.5). It is worth emphasizing that problem (1.4) can be solved for arbitrary initial values $\left(u_{0}, \eta_{0}\right)$ in a proper phase space $L^{2}(\Omega) \times L_{\mu}^{2}\left(\mathbb{R}^{+} ; H_{0}^{1}(\Omega)\right.$ ) (see Section 2 for more details), i.e., the second component $\eta_{0}$ does not necessarily depend on $u_{0}(\cdot)$. This permits us to construct a dynamical system in this phase space and prove the existence of global attractors. However, the transformed equation (1.4) is a generalization of problem (1.2), and therefore, not every solution to equation (1.4) possesses a corresponding one to (1.2). Notice that both problems are equivalent if and only if the initial value $\eta_{0}$ belongs to a proper subspace of $L_{\mu}^{2}\left(\mathbb{R}^{+} ; H_{0}^{1}(\Omega)\right)$, which coincides with the domain of the distributional derivative with respecto to $s$, denoted by $D(\mathbf{T})$ (for more details, see [10). Hence, it is natural to construct a dynamical system generated by (1.4) in the phase space $L^{2}(\Omega) \times D(\mathbf{T})$ to prove the existence of attractors to the original problem, via the above relationship (see [12, 6, 10]). Nevertheless, as far as we know, it is not possible to prove the existence of attractors in this space unless solutions are proved to have more regularity. Thus, in principle, we cannot transfer the existence of attractors for system (1.4) to the original problem (1.2).

The idea of the second alternative comes from a simple case, which was successfully applied in [1] when the kernel is $k(t)=e^{-d_{0} t}, d_{0}>0$ (non-singular kernel). Using this method, it is proved that the problem in [1 generates a dynamical system in the phase space $L_{H_{0}^{1}}^{2}$ given by the measurable functions $\varphi:(-\infty, 0] \rightarrow H_{0}^{1}(\Omega)$, such that $\int_{-\infty}^{0} e^{\gamma s}\|\varphi(s)\|_{H_{0}^{1}}^{2} d s<+\infty$, for certain $\gamma>0$. Under the construction of this phase space, there exists a global attractor to this problem (in fact, the problem in [1] is non-autonomous and the attractor is of pullback type). Notice that, for this kind of delay problems, in which
the initial value at zero may not be related to the values for negative times, the standard and more appropriate phase space to construct a dynamical system is the cartesian product $L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$ (see [2] for more details). In such a way, for any initial values $u_{0} \in L^{2}(\Omega)$ and $\varphi \in L_{H_{0}^{1}}^{2}$, there exists a unique solution to the following problem (we set $\tau=0$ since the problem is autonomous),

$$
\begin{cases}c_{0} \frac{\partial u}{\partial t}-k_{0} \Delta u-\int_{-\infty}^{t} k(t-s) \Delta u(s) d s+f(u)=g, & \text { in } \Omega \times(0, \infty)  \tag{1.6}\\ u(x, t)=0, & \text { on } \partial \Omega \times \mathbb{R} \\ u(x, 0)=u_{0}(x), & \text { in } \Omega, \\ u(x, t)=\varphi(x, t), & \text { in } \Omega \times(-\infty, 0)\end{cases}
$$

According to the regularity of solutions to the above equation, one can define a dynamical system $S(t): L^{2}(\Omega) \times L_{H_{0}^{1}}^{2} \rightarrow L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$ by the relation

$$
S(t)\left(u_{0}, \varphi\right):=\left(u\left(t ; 0, u_{0}, \varphi\right), u_{t}\left(\cdot ; 0, u_{0}, \varphi\right)\right),
$$

where $u\left(\cdot ; 0, u_{0}, \varphi\right)$ denotes the solution of problem (1.6) (see [2] for more details on this set-up). We emphasize that the two components of the dynamical system are the current state of the solution and the past history up to present, respectively, what is more sensible in a problem with delays or memory. By using this framework, the method in 11 can be successfully applied to prove the existence of attractors to problem (1.6) when $k$ is of exponential type. However, this exponential behavior may be a big restriction on the kernel $k$, consequently, on the function $\mu$, since in many real situations the latter often has singularities, for instance $k(t)=e^{-d_{0} t} t^{-\alpha}, \alpha \in(0,1)$. Therefore, it is interesting to design a technique which allows us to handle the cases with this kind of singular kernels within the context of the phase space $L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$. We will obtain this result as a consequence of the analysis performed in this paper even for the more general case of non-local problems as described below.

Let us recall now that amongst many interesting results concerning non-local differential equations, we mention the pioneering work [9, in which a model of single-species dynamics incorporating non-local effects was analyzed, comparing with the standard approach to model a single-species domain $\Omega$ of "Kolmogorov" type,

$$
\partial_{t} u=\Delta u+\lambda u g(u), \quad \text { in } \quad \Omega, \quad t>0 .
$$

Taking into account the following two natural assumptions: (i) a population in which individuals compete for a shared rapidly equilibrate resource; (ii) a population in which individuals communicate either visually or by chemical means, then the most straightforward way of introducing non-local effects is to consider, instead of $g(u)$, a "crowding" effect of the form $g(u, \bar{u})$, where

$$
\bar{u}(x, t)=\int_{\Omega} G(x, y) u(y, t) d y
$$

and $G(x, y)$ is some reasonable kernel. Reasoning in a heuristic way, Chipot et al. [5] studied the behavior of a population of bacteria with non-local term $a\left(\int_{\Omega} u\right)$ in a container. Later, Chipot et al. (cf. [3, [4) extended this term to a general non-local operator $a(l(u))$, where $l \in \mathcal{L}\left(L^{2}(\Omega) ; \mathbb{R}\right)$, for instance, if $g \in L^{2}(\Omega)$,

$$
l(u)=l_{g}(u)=\int_{\Omega} g(x) u(x) d x .
$$

Motivated by these works, the dynamics of the following non-autonomous non-local partial differential equations with delay and memory was investigated in [20] by using the Galerkin method and energy estimations,

$$
\begin{cases}\frac{\partial u}{\partial t}-a(l(u)) \Delta u=f(u)+h\left(t, u_{t}\right), & \text { in } \Omega \times(\tau, \infty)  \tag{1.7}\\ u=0, & \text { on } \partial \Omega \times \mathbb{R}, \\ u_{\tau}(x, t)=\varphi(x, t), & \text { in } \Omega \times(-\rho, 0]\end{cases}
$$

where $\Omega \subset \mathbb{R}^{N}$ is a bounded open set, $\tau \in \mathbb{R}$, the function $a \in C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$is locally Lipschitz, $f \in C(\mathbb{R}), h$ contains hereditary characteristics involving delays, and $u_{t}:(-\infty, 0] \rightarrow$ $\mathbb{R}$ is a segment of the solution given by $u_{t}(x, s)=u(x, t+s), s \leq 0$, which essentially represents the history of the solution up to time $t$. Moreover, $0<\rho \leq \infty$, which implies, the authors considered both cases, bounded and unbounded delays, for this model. However, the technique applied in [20] is the same used in [1] and, therefore, it is valid only for nonsingular memory terms of exponential kind (e.g., $k(t)=k_{1} e^{-d_{0} t}, k_{1} \in \mathbb{R}, d_{0}>0$ ), for more details, see [1]. Whereas, this technique fails to deal with various important models with memory, whose kernels have singularities.

Consequently, very recently, a new model has been considered related to long time memory differential equations containing non-local diffusion,

$$
\begin{cases}\frac{\partial u}{\partial t}-a(l(u)) \Delta u-\int_{-\infty}^{t} k(t-s) \Delta u(s) d s+f(u)=g, & \text { in } \Omega \times(\tau, \infty)  \tag{1.8}\\ u(x, t)=0, & \text { on } \partial \Omega \times \mathbb{R} \\ u(t+\tau)=\varphi(t), & \text { in } \Omega \times(-\infty, 0]\end{cases}
$$

where $\Omega \subset \mathbb{R}^{N}$ is a bounded domain with regular boundary, the function $a \in C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$ satisfies

$$
\begin{equation*}
0<m \leq a(r), \quad \forall r \in \mathbb{R} \tag{1.9}
\end{equation*}
$$

$k: \mathbb{R}^{+} \rightarrow \mathbb{R}$ is the memory kernel, with or without singularities, whose properties will be specified later, $g \in L^{2}(\Omega)$ which is independent of time. Notice that, thanks to a change of variable, the long time memory term in problem (1.8) can be interpreted as an infinite delay term,

$$
\begin{equation*}
h\left(u_{t}\right):=\int_{-\infty}^{0} k(-s) \Delta u_{t}(x, s) d s=\int_{-\infty}^{0} k(-s) \Delta u(x, t+s) d s=\int_{-\infty}^{t} k(t-s) \Delta u(x, s) d s \tag{1.10}
\end{equation*}
$$

Obviously, our model is an autonomous non-local partial differential equation. The authors first proved in [21] the existence and uniqueness of solutions to 1.8 by using the Dafermos transformation. Next, they constructed an autonomous dynamical system in the phase space $L^{2}(\Omega) \times L_{\mu}^{2}\left(\mathbb{R}^{+} ; H_{0}^{1}(\Omega)\right)$ and proved the existence of a global attractor in this space. As in the local heat equation case mentioned above, the same lack of enough regularity does not allow us to obtain an appropriate attractor for the original problem 1.8 in the phase space $L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$. Therefore, our objective is to overcome this difficulty and we succeeded by proceeding in the following way: Consider problem 1.8 with initial values $u(\tau)=u_{0}$ and $u(t+\tau)=\varphi(t)$ for $t<0$, where $\left(u_{0}, \varphi\right) \in L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$. Thus, for those kernels $\mu(\cdot)$ which guarantee that, when $\varphi \in L_{H_{0}^{1}}^{2}$ the corresponding $\eta_{\varphi}$, defined by $\eta_{\varphi}(s)=\int_{-s}^{0} \varphi(r) d r, \quad(s>0)$ belongs to the space $L_{\mu}^{2}\left(\mathbb{R}^{+} ; H_{0}^{1}(\Omega)\right)$, we can perform the Dafermos transformation and obtain the initial value problem which was already analyzed in [21], and consequently we have the existence, uniqueness and regularity of solutions in a straightforward way. Thanks to this result, we are able to construct the dynamical system in the phase space $L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$ with the help of some additional technical results. The existence of global attractor is then proved by first showing the existence of a bounded absorbing set and the proof of the asymptotic compactness property which requires an appropriate adaptation of the technique used in [1]. These results proved in the non-local problem (1.8) improve and complete the ones in [1] by simply assuming that $a(\cdot)$ is a constant, and also improve the previous literature on the local case (see, e.g., [10, 11, 12]), where it is only provided the existence of attractors for the transformed equation (1.4) but not for the original one (1.1).

The content of this paper is as follows: In Section 2, we recall some preliminaries, notations and the framework in which we will carry out our analysis. Section 3 is devoted to proving the main results of our paper. First, we state the existence and uniqueness of solutions of our problem by rewriting it as an equivalent one thanks to the Dafermos transformation. The transformed problem has already been analyzed in [21], whence our result follows immediately. However, as some estimations we need for the subsequent results are based on the ones in the proof of this existence theorem, we have included the complete proof in the Appendix (at the end of the paper). Next, we prove that our model generates an autonomous dynamical system in the phase space $L^{2}(\Omega) \times L_{H_{0}^{1}}^{2}$. Eventually, the existence of a global attractor for the dynamical system is proved by working directly on our model with memory, instead of using any result already proved in [21] for the transformed problem.

## 2. Well-posedness to a non-local differential equation with memory

The following non-local differential equation associated with singular memory will be investigated,

$$
\begin{cases}\frac{\partial u}{\partial t}-a(l(u)) \Delta u-\int_{-\infty}^{t} k(t-s) \Delta u(x, s) d s+f(u)=g(x, t), & \text { in } \Omega \times(\tau, \infty)  \tag{2.1}\\ u(x, t)=0, & \text { on } \partial \Omega \times \mathbb{R} \\ u(x, 0)=u_{0}(x), & \text { in } \Omega \\ u(x, t+\tau)=\phi(x, t), & \text { in } \Omega \times(-\infty, 0]\end{cases}
$$

where $\Omega \subset \mathbb{R}^{N}$ is a fixed bounded domain with regular boundary. The function $a \in$ $C\left(\mathbb{R} ; \mathbb{R}^{+}\right)$satisfies

$$
\begin{equation*}
0<m \leq a(r), \quad \forall r \in \mathbb{R} \tag{2.2}
\end{equation*}
$$

$k: \mathbb{R}^{+}=(0,+\infty) \rightarrow \mathbb{R}$ is the memory kernel, whose properties will be specified later. The initial values are $u_{0} \in L^{2}(\Omega)$ and $\phi \in L_{V}^{2}$ (see Section 2.2 below).

Let us define the new variables

$$
u^{t}(x, s)=u(x, t-s), \quad s \geq 0
$$

and

$$
\begin{equation*}
\eta^{t}(x, s)=\int_{0}^{s} u^{t}(x, r) d r=\int_{t-s}^{t} u(x, r) d r, \quad s \geq 0 . \tag{2.3}
\end{equation*}
$$

Assuming $k(\infty)=0$, a change of variable and a formal integration by parts yield

$$
\int_{-\infty}^{t} k(t-s) \Delta u(s) d s=-\int_{0}^{\infty} k^{\prime}(s) \Delta \eta^{t}(s) d s
$$

here and in the sequel, the prime denotes derivation with respect to variable $s$. Setting

$$
\begin{equation*}
\mu(s)=-k^{\prime}(s) \tag{2.4}
\end{equation*}
$$

the above choice of variable leads to the following non-delay system,

$$
\begin{cases}\frac{\partial u}{\partial t}-a(l(u)) \Delta u-\int_{0}^{\infty} \mu(s) \Delta \eta^{t}(s) d s+f(u)=g(x, t), & \text { in } \Omega \times(\tau, \infty)  \tag{2.5}\\ \frac{\partial}{\partial t} \eta^{t}(s)=u-\frac{\partial}{\partial s} \eta^{t}(s), & \text { in } \Omega \times(\tau, \infty) \times \mathbb{R}^{+} \\ u(x, t)=\eta^{t}(x, s)=0, & \text { on } \partial \Omega \times \mathbb{R} \times \mathbb{R}^{+} \\ u(x, \tau)=u_{0}(x), & \text { in } \Omega \\ \eta^{\tau}(x, s)=\eta_{0}(x, s), & \text { in } \Omega \times \mathbb{R}^{+}\end{cases}
$$

where, by the definition of $\eta^{t}(x, s)$ (see 2.3), it obviously follows

$$
\begin{equation*}
\eta^{\tau}(x, s)=\int_{\tau-s}^{\tau} u(x, r) d r=\int_{-s}^{0} \phi(x, r) d r:=\eta_{0}(x, s) \tag{2.6}
\end{equation*}
$$

which is the initial integrated past history of $u$ with vanishing boundary.
It is worth emphasizing that we will consider solutions of our problems in the weak (variational) sense.

### 2.1. Assumptions

In our analysis, we shall suppose the nonlinear term $f: \mathbb{R} \rightarrow \mathbb{R}$ is a polynomial of odd degree with positive leading coefficient,

$$
\begin{equation*}
f(u)=\sum_{k=1}^{2 p} f_{2 p-k} u^{k-1}, \quad p \in \mathbb{N} . \tag{2.7}
\end{equation*}
$$

This situation can be extended, without any additional difficulties, to a more general function satisfying suitable assumptions (see, for instance, [12]).

In view of the evolution problem (2.5), the variable $\mu$ is required to verify the following hypotheses:
$\left(h_{1}\right) \mu \in C^{1}\left(\mathbb{R}^{+}\right) \cap L^{1}\left(\mathbb{R}^{+}\right), \quad \mu(s) \geq 0, \quad \mu^{\prime}(s) \leq 0, \quad \forall s \in \mathbb{R}^{+} ;$
$\left(h_{2}\right) \mu^{\prime}(s)+\delta \mu(s) \leq 0, \quad \forall s \in \mathbb{R}^{+}, \quad$ for some $\delta>0$.
Remark 2.1. 1. It is straightforward to check that conditions $\left(h_{1}\right)-\left(h_{2}\right)$ are fulfilled by singular kernels given by

$$
\mu(t)=e^{-\delta t} t^{-\alpha}, t>0,
$$

for $\delta>0$ and $\alpha \in(0,1)$.
2. Restriction ( $h_{1}$ ) is equivalent to assuming $k(\cdot)$ is a non-negative, non-increasing, bounded, convex function of class $C^{2}$ vanishing at infinity. Moreover, from ( $h_{1}$ ) it easily follows that

$$
k(0)=\int_{0}^{\infty} \mu(s) d s \quad \text { is finite and non-negative. }
$$

3. Assumption $\left(h_{2}\right)$ implies that $\mu(s)$ decays exponentially. Also, this condition allows the memory kernel $k(\cdot)$ to have a singularity at $t=0$, which coincides with the intention to study problem (2.5).

### 2.2. Notations

Let $\Omega$ be a fixed bounded domain in $\mathbb{R}^{N}$. On this set, we introduce the Lebesgue space $L^{p}(\Omega)$, where $1 \leq p \leq \infty$. Besides, $W^{1, p}(\Omega)$ is the subspace of $L^{p}(\Omega)$ consisting of functions such that the first order weak derivative belongs to $L^{p}(\Omega)$. In this paper, $L^{2}(\Omega)$ is denoted by $H, H_{0}^{1}(\Omega)$ is denoted by $V$ and $H^{-1}(\Omega)$ is denoted by $V^{*}$. The norms in $H, V$ and $V^{*}$ will be denoted by $|\cdot|,\|\cdot\|$ and $\|\cdot\|_{*}$, respectively.

In view of system (2.5) and $\left(h_{1}\right)$, we need to introduce some additional notations before proving our main theorems. Let $L_{\mu}^{2}\left(\mathbb{R}^{+} ; H\right)$ be a Hilbert space of functions $w: \mathbb{R}^{+} \rightarrow H$ endowed with the inner product,

$$
\left(w_{1}, w_{2}\right)_{\mu}=\int_{0}^{\infty} \mu(s)\left(w_{1}(s), w_{2}(s)\right) d s
$$

and let $|\cdot|_{\mu}$ denote the corresponding norm. In a similar way, we introduce the inner products $((\cdot, \cdot))_{\mu},(((\cdot, \cdot)))_{\mu}$ and relative norms $\|\cdot\|_{\mu},\| \| \cdot\| \|_{\mu}$ on $L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right), L_{\mu}^{2}\left(\mathbb{R}^{+} ; V \cap H^{2}(\Omega)\right)$ respectively. It follows then that

$$
((\cdot, \cdot))_{\mu}=(\nabla \cdot, \nabla \cdot)_{\mu}, \quad \text { and } \quad(((\cdot, \cdot)))_{\mu}=(\Delta \cdot, \Delta \cdot)_{\mu} .
$$

We also define the Hilbert spaces

$$
\mathcal{H}=H \times L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)
$$

and

$$
\mathcal{V}=V \times L_{\mu}^{2}\left(\mathbb{R}^{+} ; V \cap H^{2}(\Omega)\right),
$$

which are respectively endowed with inner products

$$
\left(w_{1}, w_{2}\right)_{\mathcal{H}}=\left(w_{1}, w_{2}\right)+\left(\left(w_{1}, w_{2}\right)\right)_{\mu},
$$

and

$$
\left(w_{1}, w_{2}\right)_{\mathcal{V}}=\left(\left(w_{1}, w_{2}\right)\right)+\left(\left(\left(w_{1}, w_{2}\right)\right)\right)_{\mu}
$$

where $w_{i} \in \mathcal{H}$ or $\mathcal{V}(i=1,2)$ and usual norms.
At last, with standard notations, $\mathcal{D}(I ; X)$ is the space of infinitely differentiable $X$ valued functions with compact support in $I \subset \mathbb{R}$, whose dual space is the distribution space on $I$ with values in $X^{*}$ (dual of $X$ ), denoted by $\mathcal{D}^{\prime}\left(I ; X^{*}\right)$. For convenience, we define $L_{V}^{2}$ the space of functions $u(\cdot)$ satisfying

$$
\int_{-\infty}^{0} e^{\gamma s}\|u(s)\|^{2} d s<\infty
$$

where $0<\gamma<\min \left\{m \lambda_{1}, \delta\right\}$ and $\delta$ comes from $\left(h_{2}\right)$.

## 3. Main results

Let us start by proving a technical result which will be crucial to our analysis. To this end, we define the linear operator $\mathcal{J}: L_{V}^{2} \rightarrow L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$ by

$$
\begin{equation*}
(\mathcal{J} \phi)(s)=\int_{-s}^{0} \phi(r) d r, \quad s \in \mathbb{R}^{+} . \tag{3.1}
\end{equation*}
$$

Then we have the following result.

$$
\begin{equation*}
\|\mathcal{J} \phi\|_{L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)}^{2} \leq K_{\mu}\|\phi\|_{L_{V}^{2}}^{2} . \tag{3.2}
\end{equation*}
$$

Proof. The linearity of $\mathcal{J}$ is obvious, we only need to prove it is well defined and bounded. Indeed, taking into account the fact that $\phi \in L_{V}^{2},\left(h_{1}\right)-\left(h_{2}\right)$ and (3.1), we have

$$
\begin{aligned}
\|\mathcal{J} \phi\|_{L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)}^{2}= & \int_{0}^{\infty} \mu(s)\left\|\int_{-s}^{0} \phi(r) d r\right\|^{2} d s \\
= & \int_{0}^{1} \mu(s)\left\|\int_{-s}^{0} \phi(r) d r\right\|^{2} d s+\int_{1}^{\infty} \mu(s)\left\|\int_{-s}^{0} \phi(r) d r\right\|^{2} d s \\
\leq & \int_{0}^{1} s \mu(s) \int_{-s}^{0}\|\phi(r)\|^{2} d r d s+\mu(1) \int_{1}^{\infty} e^{-\delta(s-1)}\left\|\int_{-s}^{0} \phi(r) d r\right\|^{2} d s \\
\leq & \int_{-1}^{0}\|\phi(r)\|^{2} \int_{-r}^{1} s \mu(s) d s d r+\mu(1) e^{\delta} \int_{0}^{\infty} e^{-\delta s} s \int_{-s}^{0}\|\phi(r)\|^{2} d r d s \\
\leq & \int_{0}^{1} s \mu(s) d s \int_{-1}^{0}\|\phi(r)\|^{2} d r+\mu(1) e^{\delta} \int_{-\infty}^{0} e^{\gamma r}\|\phi(r)\|^{2} \int_{-r}^{\infty} s e^{-\gamma r} e^{-\delta s} d s d r \\
\leq & \int_{0}^{1} \mu(s) d s \int_{-1}^{0} e^{-\gamma r} e^{\gamma r}\|\phi(r)\|^{2} d r \\
& +\mu(1) e^{\delta} \int_{-\infty}^{0} e^{\gamma r}\|\phi(r)\|^{2} \int_{-r}^{\infty} s e^{\gamma s} e^{-\delta s} d s d r \\
\leq & \left(e^{\gamma} \int_{0}^{1} \mu(s) d s+\mu(1) e^{\delta}(\gamma-\delta)^{-2}\right)\|\phi\|_{L_{V}^{2}}^{2} .
\end{aligned}
$$

Lemma 3.1. Assume $\left(h_{1}\right)-\left(h_{2}\right)$ hold. Then, the operator $\mathcal{J}$ defined by (3.1) is a linear and continuous mapping. In particular, there exists a positive constant $K_{\mu}$ such that, for any $\phi \in L_{V}^{2}$, it holds


Denoting $K_{\mu}=e^{\gamma} \int_{0}^{1} \mu(s) d s+\mu(1) e^{\delta}(\gamma-\delta)^{-2}$, the proof is finished.
Remark 3.2. Notice that when we fix an initial value $\phi \in L_{V}^{2}$ for problem (2.1), then the corresponding initial value for the second component of problem 2.5) becomes $\eta_{0}:=\mathcal{J} \phi$, which belongs to $L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$ thanks to Lemma 3.1.

Before stating the existence, uniqueness and regularity of solution to our problem (2.1), we first recall a general result proved in [21] for problem (2.5) with general initial data in $H \times L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$. Let us denote

$$
z(t)=\left(u(t), \eta^{t}\right) \quad \text { and } \quad z_{0}=\left(u_{0}, \eta_{0}\right)
$$

Set

$$
\mathcal{L} z=\left(a(l(u)) \Delta u+\int_{0}^{\infty} \mu(s) \Delta \eta(s) d s, u-\eta_{s}\right)
$$

and

$$
\mathcal{G}(z)=(-f(u)+g, \quad 0) .
$$

Then problem (2.5) can be written in the following compact form,

$$
\begin{cases}z_{t}=\mathcal{L} z+\mathcal{G}(z), & \text { in } \Omega \times(\tau, \infty)  \tag{3.3}\\ z(x, t)=0, & \text { on } \partial \Omega \times(\tau, \infty) \\ z(x, \tau)=z_{0}, & \text { in } \Omega\end{cases}
$$

Now we have the following result.
Theorem 3.3 ([21]). Suppose (2.2), (2.7) and $\left(h_{1}\right)-\left(h_{2}\right)$ hold true, also let $g \in H$. In addition, assume that $a(\cdot)$ is locally Lipschitz, and there exists a positive constant $\tilde{m}$ such that,

$$
\begin{equation*}
a(s) \leq \tilde{m}, \quad \forall s \in \mathbb{R} \tag{3.4}
\end{equation*}
$$

Then:
(i) For any $z_{0} \in \mathcal{H}$, there exists a unique solution $z(\cdot)=(u(\cdot), \eta \cdot)$ to problem (3.3) which satisfies

$$
\begin{aligned}
& u(\cdot) \in L^{\infty}(\tau, T ; H) \cap L^{2}(\tau, T ; V) \cap L^{2 p}\left(\tau, T ; L^{2 p}(\Omega)\right), \quad \forall T>\tau, \\
& \eta^{\dot{\prime} \in L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right), \quad \forall T>\tau .}
\end{aligned}
$$

Furthermore, $z(\cdot) \in C(\tau, T ; \mathcal{H})$ for every $T>\tau$, and the mapping $F: z_{0} \in \mathcal{H} \rightarrow$ $z(t) \in \mathcal{H}$ is continuous for every $t \in[\tau, T]$.
(ii) For any $z_{0} \in \mathcal{V}$, the unique solution $z(\cdot)=\left(u(\cdot), \eta^{\cdot}\right)$ to problem (3.3) satisfies

$$
\begin{aligned}
& u(\cdot) \in L^{\infty}(\tau, T ; V) \cap L^{2}\left(\tau, T ; V \cap H^{2}(\Omega)\right), \quad \forall T>\tau, \\
& \eta^{\prime} \in L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V \cap H^{2}(\Omega)\right)\right), \quad \forall T>\tau .
\end{aligned}
$$

In addtion, $z(\cdot) \in C(\tau, T ; \mathcal{V})$ for every $T>\tau$.
Based on the previous theorem, we can state now the corresponding result for our original problem (2.1).

Theorem 3.4. Assume (2.2), (2.7), and $\left(h_{1}\right)-\left(h_{2}\right)$ hold. Let a $(\cdot)$ be locally Lipschitz satisfying (3.4),

$$
g \in H, \quad u_{0} \in H \quad \text { and } \quad \phi \in L_{V}^{2} .
$$

Then, there exists a unique function $z(\cdot)=\left(u(\cdot), \eta \eta^{\cdot}\right)$ satisfying

$$
\begin{aligned}
& u(\cdot) \in L^{\infty}(\tau, T ; H) \cap L^{2}(\tau, T ; V) \cap L^{2 p}\left(\tau, T ; L^{2 p}(\Omega)\right), \quad \forall T>\tau, \\
& \eta^{\prime} \in L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right), \quad \forall T>\tau,
\end{aligned}
$$

such that

$$
\partial_{t} z=\mathcal{L} z+\mathcal{G}(z)
$$

in the weak sense, and

$$
\left.z\right|_{t=\tau}=\left(u_{0}, \mathcal{J} \phi\right) .
$$

Furthermore, for every $t \in[\tau, T]$,

$$
z(t): \mathcal{H} \rightarrow \mathcal{H} \text { is a continuous mapping. }
$$

If we also assume that $u_{0} \in V, \phi \in L_{V \cap H^{2}(\Omega)}^{2}$, then

$$
\begin{aligned}
& u \in L^{\infty}(\tau, T ; V) \cap L^{2}\left(\tau, T ; V \cap H^{2}(\Omega)\right), \quad \forall T>\tau, \\
& \eta^{\prime} \in L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V \cap H^{2}(\Omega)\right)\right), \quad \forall T>\tau,
\end{aligned}
$$

and for each $t \in[\tau, T]$,

$$
z(t): \mathcal{V} \rightarrow \mathcal{V} \text { is a continuous mapping. }
$$

Proof. Thanks to Lemma 3.1, we obtain $\mathcal{J} \phi \in L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$ since $\phi \in L_{V}^{2}$. Therefore, the first statement of Theorem 3.4 holds by applying (i) in Theorem 3.3 with initial value $z_{0}=\left(u_{0}, \mathcal{J} \phi\right)$. If, in addition, we assume that initial values $u_{0} \in V$ and $\phi \in L_{V \cap H^{2}(\Omega)}^{2}$, then it is straightforward to prove that $z_{0}=\left(u_{0}, \mathcal{J} \phi\right) \in \mathcal{V}$ and the regularity result follows from statement (ii) in Theorem 3.3.

Remark 3.5. Although the proof of Theorem [3.4 follows directly from Theorem [3.3, some computations, that we need in the sequel, are based on some estimations carried out in the proof. For this reason, we have included the complete proof of Theorem 3.4 as an Appendix, so that the paper is self-contained and easier to read.

In what follows, we construct the dynamical system generated by (2.1) assuming that $g$ does not depend on $t$, which makes our problem be autonomous. Thus, the theory of autonomous dynamical systems is appropriate to carry out the analysis of the global asymptotic behavior. We emphasize that the non-autonomous case can also be studied by exploiting the theory of non-autonomous dynamical systems (either the theory of pullback attractors or the uniform attractors one). The autonomous framework is concerned with the phase space

$$
X=H \times L_{V}^{2}
$$

endowed with the norm

$$
\left\|\left(w_{1}, w_{2}\right)\right\|_{X}^{2}=\left|w_{1}\right|^{2}+\left\|w_{2}\right\|_{L_{V}^{2}}^{2}
$$

Then, thanks to Theorem 3.4, we can define a semigroup $S: \mathbb{R}^{+} \times X \rightarrow X$ by

$$
S(t)\left(u_{0}, \phi\right)=\left(u\left(t ; 0,\left(u_{0}, \mathcal{J} \phi\right)\right), u_{t}\left(\cdot ; 0,\left(u_{0}, \mathcal{J} \phi\right)\right)\right),
$$

where $\left(u\left(\cdot ; 0,\left(u_{0}, \mathcal{J} \phi\right)\right), \eta^{\prime}\right)$ is the unique solution to problem (2.5) with $u(0)=u_{0}, \eta_{0}=\mathcal{J} \phi$.
Let us first prove that the dynamical system $S$ is well defined. In what follows, we will take $\tau=0$ since we are working on autonomous dynamical system.

Lemma 3.6. Under assumptions of Theorem 3.4. if $\left(u_{0}, \phi\right) \in X$, then $S(t)\left(u_{0}, \phi\right) \in X$.
Proof. Let $\left(u_{0}, \phi\right) \in X$ and, for simplicity, denote by $\left(u(\cdot), \eta^{\cdot}\right)$ the solution to problem (2.5) corresponding to the initial value $\left(u_{0}, \mathcal{J} \phi\right)$. It follows from Theorem 3.4 that the first component $u(t)$ belongs to $H$, thus it only remains to show that the segment of solution $u_{t}(\cdot)$ belongs to $L_{V}^{2}$. Indeed,

$$
\begin{aligned}
\int_{-\infty}^{0} e^{\gamma s}\left\|u_{t}(s)\right\|^{2} d s & =\int_{-\infty}^{0} e^{\gamma s}\|u(t+s)\|^{2} d s \\
& =\int_{-\infty}^{t} e^{\gamma(\sigma-t)}\|u(\sigma)\|^{2} d \sigma \\
& =e^{-\gamma t} \int_{-\infty}^{t} e^{\gamma \sigma}\|u(\sigma)\|^{2} d \sigma \\
& =e^{-\gamma t} \int_{-\infty}^{0} e^{\gamma \sigma}\|\phi(\sigma)\|^{2} d \sigma+\int_{0}^{t} e^{\gamma(\sigma-t)}\|u(\sigma)\|^{2} d \sigma \\
& <+\infty,
\end{aligned}
$$

where the above estimation holds true since $\phi \in L_{V}^{2}$ and $u \in L^{2}(0, T ; V)$ for all $T>0$. The proof of this lemma is complete.

Lemma 3.7. Under assumptions of Theorem 3.4, there exist two positive constants $K_{1}$ and $K_{2}$, such that

$$
\begin{equation*}
\left\|S(t)\left(u_{0}, \phi\right)\right\|_{X}^{2} \leq K_{1}\left\|\left(u_{0}, \phi\right)\right\|_{X}^{2} e^{-\gamma t}+K_{2}, \quad \forall t \geq 0,\left(u_{0}, \phi\right) \in X \tag{3.5}
\end{equation*}
$$

Proof. Let $\left(u_{0}, \phi\right) \in X$ and denote by $z(\cdot)=\left(u(\cdot), \eta^{\cdot}\right)$ the solution to 2.5) corresponding to the initial value $\left(u_{0}, \mathcal{J} \phi\right)$. Now, we multiply the first equation in 2.5 ) by $u(t)$ in $H$ and the second equation in (2.5) by $\eta^{t}$ in $L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$. Then, by the same energy estimations as in the proof of Theorem 3.4 (see Appendix (3.29), we obtain

$$
\begin{aligned}
& \frac{d}{d t}\|z\|_{\mathcal{H}}^{2}+m \lambda_{1}|u|^{2}+m\|u\|^{2}+f_{0}|u|_{2 p}^{2 p}+2\left(\left(\left(\eta^{t}\right)^{\prime}, \eta^{t}\right)\right)_{\mu} \\
& \quad \leq 2 a_{0}|\Omega|+\frac{2}{\sqrt{\lambda_{1}}}|g|\|u\| \\
& \quad \leq 2 a_{0}|\Omega|+\frac{2}{m \lambda_{1}}|g|^{2}+\frac{m}{2}\|u\|^{2} .
\end{aligned}
$$

Since

$$
\begin{equation*}
2\left(\left(\left(\eta^{t}\right)^{\prime}, \eta^{t}\right)\right)_{\mu}=-\int_{0}^{\infty} \mu^{\prime}(s)\left|\nabla \eta^{t}(s)\right|^{2} d s \geq \delta \int_{0}^{\infty} \mu(s)\left|\nabla \eta^{t}(s)\right|^{2} d s, \tag{3.6}
\end{equation*}
$$

$$
\begin{equation*}
\frac{d}{d t}\|z\|_{\mathcal{H}}^{2}+\gamma\|z\|_{\mathcal{H}}^{2}+\frac{m}{2}\|u\|^{2}+f_{0}|u|_{2 p}^{2 p} \leq K_{0} \tag{3.7}
\end{equation*}
$$

where $K_{0}=2 a_{0}|\Omega|+\frac{2}{m \lambda_{1}}|g|^{2}$ and we recall that $\gamma<\min \left\{m \lambda_{1}, \delta\right\}$. Notice that inequality (3.6) has been deduced formally but can be fully justified by using mollifiers (see [12, p. $348]$ ). Now multiplying the above inequality by $e^{\gamma t}$ and integrating over $(0, t)$, neglecting the last term of the left hand side of $(3.7)$, we obtain

$$
\begin{align*}
& \|z(t)\|_{\mathcal{H}}^{2}+\frac{m}{2} \int_{0}^{t} e^{-\gamma(t-s)}\|u(s)\|^{2} d s \\
\leq & \|z(t)\|_{\mathcal{H}}^{2}+\frac{m}{2} \int_{-t}^{0} e^{\gamma s}\left\|u_{t}(s)\right\|^{2} d s \\
\leq & \left\|z_{0}\right\|_{\mathcal{H}}^{2} e^{-\gamma t}+\frac{K_{0}}{\gamma} \tag{3.8}
\end{align*}
$$

Then

$$
\begin{aligned}
\frac{m}{2}\left\|u_{t}\right\|_{L_{V}^{2}}^{2} & =\frac{m}{2} \int_{-\infty}^{0} e^{-\gamma(t-s)}\|\phi(s)\|^{2} d s+\frac{m}{2} \int_{0}^{t} e^{-\gamma(t-s)}\|u(s)\|^{2} d s \\
& \leq \frac{m}{2} e^{-\gamma t}\|\phi\|_{L_{V}^{2}}^{2}+\left\|\left(u_{0}, \mathcal{J} \phi\right)\right\|_{\mathcal{H}}^{2} e^{-\gamma t}+\frac{K_{0}}{\gamma}
\end{aligned}
$$

In view of Lemma 3.1, we have that

$$
\begin{equation*}
\left\|z_{0}\right\|_{\mathcal{H}}^{2} \leq\left|u_{0}\right|^{2}+\|\mathcal{J} \phi\|_{L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)}^{2} \leq\left|u_{0}\right|^{2}+K_{\mu}\|\phi\|_{L_{V}^{2}}^{2} \tag{3.9}
\end{equation*}
$$

Hence, (3.8)-(3.9) imply the existence of positive constants $K_{1}$ and $K_{2}$, such that

$$
\left\|S(t)\left(u_{0}, \phi\right)\right\|_{X}^{2}:=|u(t)|^{2}+\left\|u_{t}\right\|_{L_{V}^{2}}^{2} \leq K_{1}\left(\left|u_{0}\right|^{2}+\|\phi\|_{L_{V}^{2}}^{2}\right) e^{-\gamma t}+K_{2}
$$

The proof of this lemma is complete.
From Lemma 3.7, we immediately have the following result.
Corollary 3.8. The ball $B_{0}=\left\{v \in X:\|v\|_{X}^{2} \leq 2 K_{2}\right\}$ is absorbing for the semigroup $S$.
Now we shall prove the asymptotic compactness of the semigroup $S$. To this end, we first state the next result.

Lemma 3.9. Assume the hypotheses in Theorem 3.4. Let $\left\{\left(u_{0}^{n}, \phi^{n}\right)\right\}$ be a sequence, such that $\left(u_{0}^{n}, \phi^{n}\right) \rightarrow\left(u_{0}, \phi\right)$ weakly in $X$ as $n \rightarrow \infty$. Then, $S(t)\left(u_{0}^{n}, \phi^{n}\right)=\left(u^{n}(t), u_{t}^{n}\right)$ fulfills:

$$
\begin{equation*}
u^{n}(\cdot) \rightarrow u(\cdot) \quad \text { in } \quad C([r, T], H) \quad \text { for all } \quad 0<r<T \tag{3.10}
\end{equation*}
$$

$$
\begin{equation*}
u^{n}(\cdot) \rightarrow u(\cdot) \quad \text { weakly in } \quad L^{2}(0, T ; V) \quad \text { for all } \quad T>0 \tag{3.11}
\end{equation*}
$$

$$
\begin{equation*}
u^{n} \rightarrow u \quad \text { in } \quad L^{2}(0, T ; H) \quad \text { for all } \quad T>0 \tag{3.12}
\end{equation*}
$$

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left\|u_{t}^{n}-u_{t}\right\|_{L_{V}^{2}}^{2} \leq K_{5} e^{-\gamma t} \limsup _{n \rightarrow \infty}\left(\left|u_{0}^{n}-u_{0}\right|^{2}+\left\|\phi^{n}-\phi\right\|_{L_{V}^{2}}^{2}\right) \quad \text { for all } \quad t \geq 0, \tag{3.13}
\end{equation*}
$$

$$
\begin{equation*}
u_{t}^{n}(\cdot) \rightarrow u_{t}(\cdot) \quad \text { in } \quad L_{V}^{2} \quad \text { for all } \quad t \geq 0 \tag{3.15}
\end{equation*}
$$

Proof. Let $T>0$ be arbitrary. In view of (3.5) and integrating in (3.7) over $(0, T)$, we deduce that $u^{n}$ is bounded in $L^{\infty}(0, T ; H), L^{2}(0, T ; V)$ and $L^{2 p}\left(0, T ; L^{2 p}(\Omega)\right), \eta_{n}^{t}$ is bounded in $L^{\infty}\left(0, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right)$. Hence, passing to a subsequence, we have

$$
\begin{align*}
& u^{n} \rightarrow u \quad \text { weak-star in } \quad L^{\infty}(0, T ; H)  \tag{3.16}\\
& u^{n} \rightarrow u \quad \text { weakly in } \quad L^{2}(0, T ; V) \\
& u^{n} \rightarrow u \quad \text { weakly in } \quad L^{2 p}\left(0, T ; L^{2 p}(\Omega)\right) \\
& \eta_{n}^{t} \rightarrow \eta^{t} \quad \text { weak-star in } \quad L^{\infty}\left(0, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right)
\end{align*}
$$

thus (3.11) holds. Also, by the same arguments in the proof of Theorem 3.4 (see Appendix), we deduce

$$
\begin{align*}
\frac{d u^{n}}{d t} & \rightarrow \frac{d u}{d t} \quad \text { weakly in } \quad L^{2}\left(0, T ; V^{*}\right)+L^{q}\left(0, T ; L^{q}(\Omega)\right),  \tag{3.17}\\
f\left(u^{n}\right) & \rightarrow \chi \quad \text { weakly in } \quad L^{q}\left(0, T ; L^{q}(\Omega)\right) .
\end{align*}
$$

In view of (3.11) and 3.17), making use of the Compactness Theorem 18 we infer that (3.12) holds true. Thus, $u^{n}(t, x) \rightarrow u(t, x), f\left(u^{n}(t, x)\right) \rightarrow f(u(t, x))$ for a.a. $(t, x) \in$ $(0, T) \times \Omega$, so Lemma 1.3 in [16] implies that $\chi=f(u)$.

By proceeding as in the proof of Theorem 3.4, we obtain that $z(\cdot)=\left(u(\cdot), \eta^{\cdot}\right)$ is a solution to problem (2.5) with initial value $z(0)=\left(u_{0}, \mathcal{J} \phi\right)$. Thanks to the uniqueness of solution, a standard argument implies that the above convergences are true for the whole sequence.

Further, we will prove (3.10). Formally, we multiply the first equation of (2.5) by $-\Delta u(t)$ in $H$, and the second equation of 2.5 by $-\eta^{t}$ in $L_{\mu}^{2}\left(\mathbb{R}^{+} ; V \cap H^{2}(\Omega)\right)$ (these calculations can be rigorously justified via Galerkin approximations). Then, arguing as in the proof of Theorem 3.4 , we obtain

$$
\begin{aligned}
& \frac{d}{d t}\|z\|_{\mathcal{V}}^{2}+2 a(l(u))|\Delta u|^{2}+2\left(\left(\left(\eta^{t},\left(\eta^{t}\right)^{\prime}\right)\right)\right)_{\mu} \\
& =2(-f(u)+g(t),-\Delta u) \\
& \leq m|\Delta u|^{2}+\frac{2}{m}|g|^{2}+\frac{2}{m} f_{2 p-1}^{2}|\Omega|+d_{0}\|u\|^{2},
\end{aligned}
$$

Hence, by (2.2), we have

$$
\begin{equation*}
\frac{d}{d t}\|z\|_{\mathcal{V}}^{2} \leq \frac{2}{m}|g|^{2}+\frac{2}{m} f_{2 p-1}^{2}|\Omega|+d_{0}\|u\|^{2} \leq K_{3}\left(1+\|u\|^{2}\right), \tag{3.18}
\end{equation*}
$$

where we have used the notation,

$$
K_{3}=\max \left\{\frac{2}{m}|g|^{2}+\frac{2}{m} f_{2 p-1}^{2}|\Omega|, \quad d_{0}\right\} .
$$

Integrating in (3.7) over $(t, t+r)$ for $t \geq 0,0<r<T-t$ and using (3.8), we deduce that

$$
\begin{equation*}
\int_{t}^{t+r}\|u\|^{2} d s \leq \frac{2 K_{0}}{m} r+\frac{2}{m}\|z(t)\|_{\mathcal{H}}^{2} \leq K_{4}(1+r), \quad \forall t \geq 0 \tag{3.19}
\end{equation*}
$$

where we have used the notation

$$
K_{4}=\max \left\{\frac{2 K_{0}}{m}, \frac{2}{m}\left\|z_{0}\right\|_{\mathcal{H}}^{2}+\frac{2 K_{0}}{m \gamma}\right\} .
$$

We integrate in (3.18) over $(s, t+r)$, where $s \in(t, t+r)$. Thus, by (3.19),

$$
\|z(t+r)\|_{\mathcal{V}}^{2} \leq\|z(s)\|_{\mathcal{V}}^{2}+K_{3} r+K_{3} K_{4}(1+r) .
$$

Integrating the above inequality now again over $(t, t+r)$ in $s$, with the help of (3.19), we have

$$
r\|z(t+r)\|_{\mathcal{V}}^{2} \leq\left\|z_{0}\right\|_{\mathcal{V}}^{2} r+2 K_{3} r^{2}+\left(K_{3}+1\right) K_{4} r(1+r), \quad \forall t \geq 0,
$$

thus, $\|z(t)\|_{\mathcal{V}}$ is uniformly bounded in $[r, T]$. We observe that by a standard argument (see [1, p.195]), for any sequence $t_{n} \rightarrow t_{0}$ as $n \rightarrow \infty, t_{n}, t_{0} \in[0, T], u^{n}\left(t_{n}\right) \rightarrow u\left(t_{0}\right)$ weakly in $V$. Then the compact embedding $V \subset H$ ensures $u^{n}\left(t_{n}\right) \rightarrow u\left(t_{0}\right)$ strongly in $H$, for all $t_{n}, t_{0} \in[r, T]$ and $t_{n} \rightarrow t_{0}$ as $n \rightarrow \infty$, therefore (3.10) holds true.

Define the functions $w^{n}=z^{n}-z, \beta_{n}^{t}=\eta_{n}^{t}-\eta^{t}$, similarly to the uniqueness step in the proof of Theorem 3.4. Step 5 in Appendix, we have

$$
\begin{align*}
& \frac{d}{d t}\left\|w^{n}\right\|_{\mathcal{H}}^{2}+2\left(\left(\left(\beta_{n}^{t}\right)^{\prime}, \beta_{n}^{t}\right)\right)_{\mu}  \tag{3.20}\\
& \leq-2 \int_{\Omega}\left(f\left(u^{n}\right)-f(u)\right)\left(u^{n}-u\right) d x-\int_{\Omega}\left(a\left(l\left(u^{n}\right)\right) \nabla u^{n}-a(l(u)) \nabla u\right) \cdot \nabla\left(u^{n}-u\right) d x .
\end{align*}
$$

Since $a$ is locally Lipschitz, by (2.2) and the Young inequality, we have

$$
\begin{align*}
& -2 \int_{\Omega}\left(a\left(l\left(u^{n}\right)\right) \nabla u^{n}-a(l(u)) \nabla u\right) \cdot \nabla\left(u^{n}-u\right) d x \\
& =-2 \int_{\Omega} a\left(l\left(u^{n}\right)\right)\left|\nabla\left(u^{n}-u\right)\right|^{2} d x-2\left(a\left(l\left(u^{n}\right)\right)-a(l(u))\right) \int_{\Omega} \nabla u \cdot \nabla\left(u^{n}-u\right) d x \\
& \leq-2 m\left\|u^{n}-u\right\|^{2}+2 L_{a}(R)|l|\left|u^{n}-u\right|\|u\|\left\|u^{n}-u\right\| \\
& \leq(\alpha-2 m)\left\|u^{n}-u\right\|^{2}+\frac{L_{a}^{2}(R)|l|^{2}}{\alpha}\left|u^{n}-u\right|^{2}\|u\|^{2} \tag{3.21}
\end{align*}
$$

where $\alpha \leq\left(m \lambda_{1}-\gamma\right) / \lambda_{1}$, and for all $n \geq 1, t \geq 0$, we choose $R>0$ such that $\left|u^{n}(t)\right|,|u(t)| \leq$ $R$ (cf. (3.10)). By (3.6), (3.20) and (3.21), we have

$$
\begin{aligned}
& \frac{d}{d t}\left\|w^{n}\right\|_{\mathcal{H}}^{2}+\gamma\left\|w^{n}\right\|_{\mathcal{H}}^{2}+m\left\|u^{n}-u\right\|^{2} \\
& \leq \frac{d}{d t}\left\|w^{n}\right\|_{\mathcal{H}}^{2}+(2 m-\alpha)\left\|u^{n}-u\right\|^{2}+\delta \int_{0}^{\infty} \mu(s)\left|\nabla \beta_{n}^{t}(s)\right|^{2} d s \\
& \leq \frac{L_{a}^{2}(R)|l|^{2}}{\alpha}\left|u^{n}-u\right|^{2}\|u\|^{2}-2 \int_{\Omega}\left(f\left(u^{n}\right)-f(u)\right)\left(u^{n}-u\right) d x
\end{aligned}
$$

where we have used that $\gamma \leq \min \left\{(m-\alpha) \lambda_{1}, \delta\right\}$ by the choice of $\alpha$. Multiplying by $e^{\gamma t}$ on both sides of the above inequality and integrating over $(0, t)$, we obtain

$$
\begin{aligned}
& \left\|w^{n}(t)\right\|_{\mathcal{H}}^{2}+m \int_{0}^{t} e^{-\gamma(t-s)}\left\|u^{n}-u\right\|^{2} d s \\
& \leq e^{-\gamma t}\left\|w^{n}(0)\right\|_{\mathcal{H}}^{2}+\frac{L_{a}^{2}(R)|l|^{2}}{\alpha} \int_{0}^{t} e^{-\gamma(t-s)}\left|u^{n}-u\right|^{2}\|u\|^{2} d s \\
& \quad-2 \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega}\left(f\left(u^{n}\right)-f(u)\right)\left(u^{n}-u\right) d x d s .
\end{aligned}
$$

On the one hand, by 3.10 , we know that $\left|u^{n}(s)-u(s)\right|^{2}\|u(s)\|^{2} \rightarrow 0$ for a.e. $s \in(0, t)$. On the other hand, $e^{-\gamma(t-s)}\left|u^{n}(s)-u(s)\right|^{2}\|u(s)\|^{2}$ is bounded by the integrable function $4 R^{2} e^{-\gamma(t-s)}\|u(s)\|^{2}$. Hence, Lebesgue's theorem implies that

$$
\int_{0}^{t} e^{-\gamma(t-s)}\left|u^{n}-u\right|^{2}\|u\|^{2} d s \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty
$$

Since $f\left(u^{n}\right) \rightarrow f(u)$ weakly in $L^{q}\left(0, T ; L^{q}(\Omega)\right)$, it follows that

$$
\int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega}\left(f\left(u^{n}\right)-f(u)\right) u d x d s \rightarrow 0 \text { as } n \rightarrow \infty
$$

Furthermore, as $f\left(u^{n}(t, x)\right) u^{n}(t, x) \geq-\kappa_{1}+\kappa_{2}\left|u^{n}(t, x)\right|^{2 p}$ (see (3.28)) and $u^{n}(t, x) \rightarrow$ $u(t, x), f\left(u^{n}(t, x)\right) \rightarrow f(u(t, x))$ for a.a. $(t, x) \in(0, T] \times \Omega$, Lebesgue-Fatous's theorem
implies

$$
\begin{aligned}
& \limsup _{n \rightarrow \infty}\left(-2 \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega} f\left(u^{n}\right) u^{n} d x d s\right) \\
& \leq-2 \liminf _{n \rightarrow \infty} \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega} f\left(u^{n}\right) u^{n} d x d s \\
& \leq-2 \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega} \liminf _{n \rightarrow \infty} f\left(u^{n}\right) u^{n} d x d s \\
& =2 \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega} f(u) u d x d s .
\end{aligned}
$$

This inequality, together with

$$
\begin{equation*}
\int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega} f(u)\left(u^{n}-u\right) d x d s \rightarrow 0 \text { as } n \rightarrow \infty \tag{3.22}
\end{equation*}
$$

231 shows that

$$
\limsup _{n \rightarrow \infty}\left(-2 \int_{0}^{t} e^{-\gamma(t-s)} \int_{\Omega}\left(f\left(u^{n}\right)-f(u)\right) u^{n} d x d s\right) \leq 0 \text { as } n \rightarrow \infty .
$$

232 Notice that (3.22) follows from the facts $f(u(\cdot)) \in L^{q}\left(0, T ; L^{q}(\Omega)\right)$ and $u^{n} \rightarrow u$ weakly in ${ }_{233} \quad L^{2 p}\left(0, T ; L^{2 p}(\Omega)\right)$.

Collecting all inequalities and using (3.2),

$$
\begin{aligned}
& \limsup _{n \rightarrow \infty} \int_{0}^{t} e^{-\gamma(t-s)}\left\|u^{n}(s)-u(s)\right\|^{2} d s \\
& \leq \frac{1}{m} e^{-\gamma t} \limsup _{n \rightarrow \infty}\left\|w^{n}(0)\right\|_{\mathcal{H}}^{2} \\
& =\frac{1}{m} e^{-\gamma t} \limsup _{n \rightarrow \infty}\left(\left|u^{n}(0)-u_{0}\right|^{2}+\int_{0}^{\infty} \mu(s)\left\|\beta_{n}^{0}(s)\right\|^{2} d s\right) \\
& \leq \frac{1}{m} e^{-\gamma t} \limsup _{n \rightarrow \infty}\left(\left|u^{n}(0)-u_{0}\right|^{2}+K_{\mu} \int_{-\infty}^{0} e^{\gamma s}\left\|\phi^{n}(s)-\phi(s)\right\|^{2} d s\right) .
\end{aligned}
$$

Finally, (3.13) follows from

$$
\begin{aligned}
\left\|u_{t}^{n}-u_{t}\right\|_{L_{V}^{2}}^{2} & =\int_{-t}^{0} e^{\gamma s}\left\|u^{n}(t+s)-u(t+s)\right\|^{2} d s+\int_{-\infty}^{-t} e^{\gamma s}\left\|u^{n}(t+s)-u(t+s)\right\|^{2} d s \\
& =\int_{0}^{t} e^{-\gamma(t-s)}\left\|u^{n}(s)-u(s)\right\|^{2} d s+e^{-\gamma t} \int_{-\infty}^{0} e^{\gamma s}\left\|\phi^{n}(s)-\phi(s)\right\|^{2} d s
\end{aligned}
$$

234 If $\left(u_{0}^{n}, \phi^{n}\right) \rightarrow\left(u_{0}, \phi\right)$ in $X$, then (3.13) implies (3.14) and (3.15).
As a consequence, we obtain the continuous dependence with respect to the initial data.

Proof. Let $B \subset X$ be a bounded set, we need to prove that for any sequences $\left\{\left(y_{n}, \phi_{n}\right)\right\}_{n \in \mathbb{N}} \subset B$ and $t_{n} \rightarrow+\infty$ as $n \rightarrow+\infty$, the sequence $\left\{S\left(t_{n}\right)\left(y_{n}, \phi_{n}\right)\right\}_{n \in \mathbb{N}}$ is relatively compact. Recall that

$$
S\left(t_{n}\right)\left(y_{n}, \phi_{n}\right)=\left(u\left(t_{n} ; 0,\left(y_{n}, \mathcal{J} \phi_{n}\right)\right), u_{t_{n}}\left(\cdot ; 0,\left(y_{n}, \mathcal{J} \phi_{n}\right)\right)\right):=\left(u^{n}\left(t_{n}\right), u_{t_{n}}^{n}(\cdot)\right)
$$

Pick now $T>0$, and assume that $t_{n}>T$ for all $n \in \mathbb{N}$ (there is no loss of generality in assuming this since $\left.t_{n} \rightarrow+\infty\right)$. Now we can define $v^{n}(t)=u^{n}\left(t+t_{n}-T\right)$, observe that $v^{n}(T)=u^{n}\left(t_{n}\right)$ and $v_{T}^{n}(t)=v^{n}(T+t)=u^{n}\left(t+t_{n}\right)=u_{t_{n}}^{n}(t)$. Therefore

$$
S\left(t_{n}\right)\left(y_{n}, \phi_{n}\right)=\left(u^{n}\left(t_{n}\right), u_{t_{n}}^{n}(\cdot)\right)=\left(v^{n}(T), v_{T}^{n}(\cdot)\right) .
$$

Let us denote now

$$
\mathcal{Y}_{n}=\left(v^{n}(T), v_{T}^{n}\right)=\left(u^{n}\left(t_{n}\right), u_{t_{n}}^{n}(\cdot)\right), \xi_{n}^{T}=\left(v^{n}(0), v_{0}^{n}(\cdot)\right)=\left(u^{n}\left(t_{n}-T\right), u_{t_{n}-T}^{n}(\cdot)\right) .
$$

By Lemma 3.7, the sequences $\left\{\mathcal{Y}_{n}\right\},\left\{\xi_{n}^{T}\right\}$ are bounded in $X$, so up to a subsequence $\mathcal{Y}_{n} \rightarrow \mathcal{Y}:=(y, \phi), \xi_{n}^{T} \rightarrow \xi^{T}$ weakly in $X$. In addition, by Lemma 3.9, $\mathcal{V}(t):=S(t) \xi^{T}=$ $\left(v(t), v_{t}(\cdot)\right)$ satisfies (3.10)-(3.13). It follows from the above convergences that, $\phi=v_{T}$ in $L_{V}^{2}$ and $y=v_{T}(0), \phi(s)=v_{T}(s)$ for almost all $s \in(-\infty, 0)$. Also, in view of 3.10) we infer that

$$
u^{n}\left(t_{n}\right)=v^{n}(T) \rightarrow v(T)=y .
$$

Hence, in order to prove that $\mathcal{Y}_{n} \rightarrow \mathcal{Y}$ in $X$, it remains to check that $u_{t_{n}}^{n}(\cdot) \rightarrow \phi$ in $L_{V}^{2}$ (up to a subsequence). Notice that $u_{t_{n}}^{n}(\cdot)=v_{T}^{n}$ for all $t_{n}>T$ and $v_{T}=\phi$. Thanks to (3.13) we have, for each $T>0$,

$$
\begin{aligned}
\limsup _{n \rightarrow \infty}\left\|u_{t_{n}}^{n}(\cdot)-\phi\right\|_{L_{V}^{2}}^{2} & =\limsup _{n \rightarrow \infty}\left\|v_{T}^{n}-v_{T}\right\|_{L_{V}^{2}}^{2} \\
& \leq K_{5} e^{-\gamma(T-\tau)} \limsup _{n \rightarrow \infty}\left(\left\|\xi_{n}^{T}-\xi^{T}\right\|_{X}^{2}\right) \\
& \leq \widetilde{K} e^{-\gamma T}
\end{aligned}
$$

where the last inequality follows from Lemma 3.7 . For every $k>0$, there exists $T:=T(k)$ such that for all $T \geq T(k)$,

$$
\limsup _{n \rightarrow \infty}\left\|u_{t_{n}}^{n}(\cdot)-\phi\right\|_{L_{V}^{2}}^{2}=\limsup _{n \rightarrow \infty}\left\|v_{T}^{n}-v_{T}\right\|_{L_{V}^{2}}^{2} \leq \frac{1}{k}
$$

Taking $k \rightarrow \infty$ and using a diagonal argument, we obtain that there exists a subsequence $\left\{u_{t_{n_{k}}}^{n_{k}}(\cdot)\right\}$ such that $u_{t_{n_{k}}}^{n_{k}}(\cdot) \rightarrow \phi$ in $L_{V}^{2}$.

By Corollaries 3.8, 3.10 and Lemma 3.11 the general theory of attractors (see [14, Theorem 3.1]) implies the following result.

Theorem 3.12. Under the assumptions of Theorem 3.4 , the semigroup $S$ possesses a global connected attractor $\mathcal{A} \subset X$.

As a straightforward consequence of the previous results, we can provide information for the local problem analyzed, amongst others, in the papers [10, 11, 12] by simply assuming that $a(\cdot)$ is a constant function.

Corollary 3.13. Under the hypotheses of Theorem 3.4, assume also that $a(t)=k_{0}>0$ for all $t \geq 0$. Then the local problem (2.1) poseesses a global connected attractor $\mathcal{A} \subset X$.

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## Appendix

Proof of Theorem 3.4. We follow a standard Faedo-Galerkin method. Recall that there exists a smooth orthonormal basis $\left\{w_{j}\right\}_{j=1}^{\infty}$ in $H$ which also belongs to $V \cap L^{2 p}(\Omega)$ ([12, Proposition 2.3]). Let us take a complete set of normalized eigenfunctions for $-\Delta$ in $V$, such that $-\Delta w_{j}=\lambda_{j} w_{j}\left(\lambda_{j}\right.$ the eigenvalue corresponding to $\left.w_{j}\right)$. Next we will select an orthonormal basis $\left\{\zeta_{j}\right\}_{j=1}^{\infty}$ of $L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$ which also belongs to $\mathcal{D}\left(\mathbb{R}^{+} ; V\right)$.

The proof is divided into 6 steps.
Step 1. (Faedo-Galerkin scheme) Fix $T>\tau$, for a given integer $n$, denote by $P_{n}$ and $Q_{n}$ the projections on the subspaces

$$
\operatorname{span}\left\{w_{1}, \cdots, w_{n}\right\} \subset V \quad \text { and } \quad \operatorname{span}\left\{\zeta_{1}, \cdots, \zeta_{n}\right\} \subset L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right),
$$

respectively. We look for a function $z_{n}=\left(u_{n}, \eta_{n}^{t}\right)$ of the form

$$
u_{n}(t)=\sum_{j=1}^{n} b_{j}(t) w_{j} \quad \text { and } \quad \eta_{n}^{t}(s)=\sum_{j=1}^{n} c_{j}(t) \zeta_{j}(s),
$$

satisfying

$$
\left\{\begin{array}{l}
\left(\partial_{t} z_{n},\left(w_{k}, \zeta_{j}\right)\right)_{\mathcal{H}}=\left(\mathcal{L} z_{n},\left(w_{k}, \zeta_{j}\right)\right)+\left(\mathcal{G}(z),\left(w_{k}, \zeta_{j}\right)\right), \quad k, j=0, \cdots, n,  \tag{3.23}\\
\left.z_{n}\right|_{t=\tau}=\left(P_{n} u_{0}, Q_{n} \eta_{0}\right)
\end{array}\right.
$$

for a.e. $\tau \leq t \leq T$, where $w_{0}$ and $\zeta_{0}$ are the zero vectors in the respective spaces. Taking $\left(w_{k}, \zeta_{0}\right)$ and $\left(w_{0}, \zeta_{k}\right)$ in (3.23), applying the divergence theorem, we derive a system of ODE in the variables

$$
\left\{\begin{array}{l}
\frac{d}{d t} b_{k}(t)=-\lambda_{k} a\left(l\left(\sum_{j=1}^{n} b_{j}(t) w_{j}\right)\right) b_{k}-\sum_{j=1}^{n} c_{j}\left(\left(\zeta_{j}, w_{k}\right)\right)_{\mu}-\left(f\left(\sum_{j=1}^{n} b_{j}(t) w_{j}\right), w_{k}\right)+\left(g, w_{k}\right),  \tag{3.24}\\
\frac{d}{d t} c_{k}(t)=\sum_{j=1}^{n} b_{j}\left(\left(w_{j}, \zeta_{k}\right)\right)_{\mu}-\sum_{j=1}^{n} c_{j}\left(\left(\zeta_{j}^{\prime}, \zeta_{k}\right)\right)_{\mu}
\end{array}\right.
$$

hence,

$$
\begin{equation*}
\left(\mathcal{G}\left(z_{n}\right), z_{n}\right)_{\mathcal{H}}=\left(-f\left(u_{n}\right)+g, u_{n}\right) \leq-\frac{1}{2} f_{0}\left|u_{n}\right|_{2 p}^{2 p}+a_{0}|\Omega|+\left(g, u_{n}\right) . \tag{3.28}
\end{equation*}
$$

It follows from $(2.2),(3.26)-(3.28)$ and the Young inequality that

$$
\begin{equation*}
\frac{d}{d t}\left\|z_{n}\right\|_{\mathcal{H}}^{2}+2 m\left|\nabla u_{n}\right|^{2}+2\left(\left(\left(\eta_{n}^{t}\right)^{\prime}, \eta_{n}^{t}\right)\right)_{\mu}+f_{0}\left|u_{n}\right|_{2 p}^{2 p} \leq 2 a_{0}|\Omega|+\frac{1}{m \lambda_{1}}|g|^{2}+m\left|\nabla u_{n}\right|^{2} . \tag{3.29}
\end{equation*}
$$

Integration by parts and $\left(h_{1}\right)$ yield that,

$$
2\left(\left(\left(\eta_{n}^{t}\right)^{\prime}, \eta_{n}^{t}\right)\right)_{\mu}=-\int_{0}^{\infty} \mu^{\prime}(s)\left|\nabla \eta_{n}^{t}(s)\right|^{2} d s \geq 0
$$

thus the third term of the right hand side of 3.29 can be neglected, we obtain

$$
\frac{d}{d t}\left\|z_{n}\right\|_{\mathcal{H}}^{2}+m\left|\nabla u_{n}\right|^{2}+f_{0}\left|u_{n}\right|_{2 p}^{2 p} \leq 2 a_{0}|\Omega|+\frac{1}{m \lambda_{1}}|g|^{2}
$$

Passing to a subsequence, there exists a function $z=(u, \eta)$ such that

$$
\left\{\begin{array}{lll}
u_{n} \rightarrow u & \text { weak-star in } & L^{\infty}(\tau, T ; H) ;  \tag{3.31}\\
u_{n} \rightarrow u & \text { weakly in } & L^{2}(\tau, T ; V) ; \\
u_{n} \rightarrow u & \text { weakly in } & L^{2 p}\left(\tau, T ; L^{2 p}(\Omega)\right) ; \\
\eta_{n}^{t} \rightarrow \eta^{t} & \text { weak-star in } & L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right)
\end{array}\right.
$$

Step 3. (Passage to limit) For a fixed integer $m$, choose a function

$$
v=(\sigma, \pi) \in \mathcal{D}\left((\tau, T) ; V \cap L^{2 p}(\Omega)\right) \times \mathcal{D}\left((\tau, T) ; \mathcal{D}\left(\mathbb{R}^{+} ; V\right)\right)
$$

of the form

$$
\sigma(t)=\sum_{j=1}^{m} \tilde{b}_{j}(t) w_{j} \quad \text { and } \quad \pi^{t}(s)=\sum_{j=1}^{m} \tilde{c}_{j}(t) \zeta_{j}(s),
$$

where $\left\{\tilde{b}_{j}\right\}_{j=1}^{m}$ and $\left\{\tilde{c}_{j}\right\}_{j=1}^{m}$ are given functions in $\mathcal{D}(\tau, T)$, then 3.23 holds with $(\sigma, \pi)$ in place of $\left(\omega_{k}, \zeta_{j}\right)$.

Our main target is to prove problem (2.5) has a solution in the weak sense, i.e., for arbitrary $v \in \mathcal{D}\left((\tau, T) ; V \cap L^{2 p}(\Omega)\right) \times \mathcal{D}\left((\tau, T) ; \mathcal{D}\left(\mathbb{R}^{+} ; V\right)\right)$ (here, specially, we pick up $v=$ $(\sigma, \pi) \in \mathcal{D}(\tau, T)$ as a test function), the following equality

$$
\begin{align*}
\int_{\tau}^{t}\left(\partial_{r} z_{n}, v\right)_{\mathcal{H}} d r=\int_{\tau}^{t}[ & -a\left(l\left(u_{n}\right)\right)\left(\nabla u_{n}, \nabla \sigma\right)-\left(\left(\eta_{n}^{t}, \sigma\right)\right)_{\mu}-\left(f\left(u_{n}\right), \sigma\right)  \tag{3.32}\\
& \left.+(g, \sigma)+\left(\left(u_{n}, \pi^{t}\right)\right)_{\mu}-\ll\left(\eta_{n}^{t}\right)^{\prime}, \pi^{t} \gg\right] d r
\end{align*}
$$

holds in the sense of $\mathcal{D}^{\prime}(\tau, T)$. Here, we denote by $\ll \cdot \cdot \gg$ the duality map between $H_{\mu}^{1}\left(\mathbb{R}^{+} ; V\right)$ and its dual space.

Firstly, using the same argument as in [20, Theorem 2.7] and (3.31)2, we know

$$
\int_{\tau}^{t} a\left(l\left(u_{n}\right)\right)\left(\nabla u_{n}, \nabla \sigma\right) d r \rightarrow \int_{\tau}^{t} a(l(u))(\nabla u, \nabla \sigma) d r \quad \text { as } \quad n \rightarrow \infty .
$$

Similarly, by means of 3.31$)_{4}$ and 3.31$)_{2}$, we have

$$
\int_{\tau}^{t}\left(\left(\eta_{n}^{t}, \sigma\right)\right)_{\mu} d r \rightarrow \int_{\tau}^{t}\left(\left(\eta^{t}, \sigma\right)\right)_{\mu} d r \quad \text { as } \quad n \rightarrow \infty
$$

and

$$
\int_{\tau}^{t}\left(\left(u_{n}, \pi^{t}\right)\right)_{\mu} d r \rightarrow \int_{\tau}^{t}\left(\left(u, \pi^{t}\right)\right)_{\mu} d r \quad \text { as } \quad n \rightarrow \infty,
$$

Notice that, for every $v \in L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)$, making use of integration by parts, we derive

$$
\begin{equation*}
\ll v^{\prime}, \pi^{t} \gg=-\int_{0}^{\infty} \mu^{\prime}(s)\left(\nabla v(s), \nabla \pi^{t}(s)\right) d s-\int_{0}^{\infty} \mu(s)\left(\nabla v(s), \nabla\left(\pi^{t}\right)^{\prime}(s)\right) d s . \tag{3.33}
\end{equation*}
$$

Replacing $v$ by $\eta_{n}^{t}$ in (3.33), together with (3.31) 4 , it is clear the right hand side of (3.33) converges to $\ll\left(\eta^{t}\right)^{\prime}, \pi^{t} \gg$ as $n \rightarrow \infty$.

Thirdly, we are going to prove that

$$
\lim _{n \rightarrow \infty} \int_{\tau}^{T} \int_{\Omega}\left|f\left(u_{n}\right) \sigma\right| d x d t=\int_{\tau}^{T} \int_{\Omega}|f(u) \sigma| d x d t
$$

Based on the dominated convergence theorem, it is sufficient to show

$$
f\left(u_{n}(t, x)\right) \rightarrow f(u(t, x)) \quad \text { for a.e. }(t, x) \in(\tau, T) \times \Omega \text {, }
$$

and

$$
\left|f\left(u_{n}\right)\right|_{L^{q}((\tau, T) \times \Omega)} \leq C,
$$

where $q=\frac{2 p}{2 p-1} \in(1,2)$, which is the conjugate exponent of $2 p$ and the constant $C$ is independent of $n$. Observe that

$$
\begin{align*}
\left\|\partial_{t} u_{n}\right\|_{L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right)} \leq & \left\|a\left(l\left(u_{n}\right)\right) \Delta u_{n}\right\|_{L^{2}\left(\tau, T ; V^{*}\right)}+\left\|\int_{0}^{\infty} \mu(s) \Delta \eta_{n}^{t}(s) d s\right\|_{L^{2}\left(\tau, T ; V^{*}\right)} \\
& +\|g\|_{V^{*}}+\|f(u)\|_{L^{q}\left(\tau, T ; L^{q}(\Omega)\right)} . \tag{3.34}
\end{align*}
$$

It follows from $(2.7)$, there exists a constant $K>0$ such that

$$
\begin{equation*}
\left|f\left(u_{n}\right)\right|^{q} \leq K\left(1+\left|u_{n}\right|^{2 p}\right) \tag{3.35}
\end{equation*}
$$

Together with (3.4), (3.31) and the assumption $g \in H$, we know that $\left\{\partial_{t} u_{n}\right\}$ is bounded in $L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right)$. Thus, up to a subsequence, we infer

$$
\begin{equation*}
\partial_{t} u_{n} \rightarrow \tilde{u} \quad \text { weakly in } \quad L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right) \tag{3.36}
\end{equation*}
$$

By a standard argument we infer that $\tilde{u}=u_{t}$. Since

$$
L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right) \subset L^{q}\left(\tau, T ; V^{*}+L^{q}(\Omega)\right)
$$

and

$$
L^{2}(\tau, T ; V) \subset L^{q}(\tau, T ; V)
$$

by (3.31) and (3.36), we deduce

$$
\begin{equation*}
u_{n} \rightarrow u \quad \text { weakly in } \quad W^{1, q}\left(\tau, T ; V^{*}+L^{q}(\Omega)\right) \cap L^{q}(\tau, T ; V) \tag{3.37}
\end{equation*}
$$

Applying a compactness argument [16], we derive that the injection

$$
W^{1, q}\left(\tau, T ; V^{*}+L^{q}(\Omega)\right) \cap L^{q}(\tau, T ; V) \hookrightarrow L^{q}\left(\tau, T ; L^{q}(\Omega)\right)
$$

is compact. Therefore, (3.37) implies that

$$
u_{n} \rightarrow u \quad \text { strongly in } \quad L^{q}\left(\tau, T ; L^{q}(\Omega)\right)
$$

By the continuity of $f$ we obtain that (up to a subsequence)

$$
f\left(u_{n}(t, x)\right) \rightarrow f(u(t, x)) \quad \text { for a.e. }(t, x) \in(\tau, T) \times \Omega
$$

In virtue of (3.35), we obtain

$$
\left|f\left(u_{n}\right)\right|_{L^{q}((\tau, T) \times \Omega)}^{q}=\int_{\tau}^{T} \int_{\Omega}\left|f\left(u_{n}\right)\right|^{q} d x d t \leq K|\Omega|(T-\tau)+K \int_{\tau}^{T}\left|u_{n}\right|_{2 p}^{2 p} d t
$$

which is bounded uniformly with respect to $n$.
Eventually, by a standard argument, we derive

$$
\partial_{t} z_{n} \rightarrow z_{t} \quad \text { in } \quad \mathcal{D}^{\prime}\left(\tau, T ; V \cap L^{2 p}(\Omega)\right) \times \mathcal{D}^{\prime}\left(\tau, T ; \mathcal{D}\left(\mathbb{R}^{+} ; V\right)\right)
$$

Step 4. (Continuity of solution) By (3.33)-3.34), it is immediate to see that $z_{t}=\left(u_{t}, \eta_{t}\right)$ fulfills

$$
\begin{aligned}
& u_{t} \in L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right) \\
& \eta_{t} \in L^{2}\left(\tau, T ; H_{\mu}^{-1}\left(\mathbb{R}^{+} ; V\right)\right)
\end{aligned}
$$

where $L^{2}\left(\tau, T ; V^{*}\right)+L^{q}\left(\tau, T ; L^{q}(\Omega)\right)$ is the dual space of $L^{2}(\tau, T, V) \cap L^{2 p}\left(\tau, T ; L^{2 p}(\Omega)\right)$. Using a slightly modified version of [19, Lemma III.1.2], together with (3.31), we infer that $u \in C([\tau, T] ; H)$.

As for the second component, by means of the same argument as [12, Theorem, Section 2], we obtain that $\eta^{t} \in C\left([\tau, T] ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; V\right)\right)$. Thus, $z(\tau)$ makes sense, and the equality $z(\tau)=z_{0}$ follows from the fact that $\left(P_{n} u_{0}, Q_{n} \eta_{0}\right)$ converges to $z_{0}$ strongly.

Step 5. (Continuity with respect to the initial value and uniqueness) Let $z_{1}=\left(u_{1}, \eta_{1}\right)$ and $z_{2}=\left(u_{2}, \eta_{2}\right)$ be the two solutions of $(3.3)$ with initial data $z_{10}$ and $z_{20}$, respectively. Due to the a priori estimates on the first component of solutions $u$, see (3.30), together with the fact that $u \in C(\tau, T ; H)$, we can ensure that there exists a bounded set $S \subset H$, such that $u_{i}(t) \in S$ for all $t \in[\tau, T]$ and $i=1,2$. In addition, taking into account that $l \in \mathcal{L}(H ; \mathbb{R})$, we have $\left\{l\left(u_{i}(t)\right)\right\}_{t \in[\tau, T]} \subset[-R, R]$ for $i=1,2$, for some $R>0$. Therefore, let $\bar{z}=z_{1}-z_{2}=(\bar{u}, \bar{\eta})=\left(u_{1}-u_{2}, \eta_{1}-\eta_{2}\right)$ and $\bar{z}_{0}=z_{10}-z_{20}$. Thanks to (2.2), the locally Lipschitz continuity of function $a$ with Lipschitz constant $L_{a}(R)$ and the Poincaré inequality, we have

$$
\begin{align*}
\frac{d}{d t}\|\bar{z}\|_{\mathcal{H}}^{2} \leq & 2 a\left(l\left(u_{1}\right)\right)|\nabla \bar{u}|^{2}+2 L_{a}(R)|l|\left|\bar{u} \| \nabla u_{2}\right||\nabla \bar{u}| \\
& -2<f\left(u_{1}\right)-f\left(u_{2}\right), \bar{u}>_{L^{p, q}}-2\left(\left((\bar{\eta})^{\prime}, \bar{\eta}\right)\right)_{\mu} \\
\leq & -2 m|\nabla \bar{u}|^{2}+2 L_{a}(R)\left|l\|\bar{u}\| \nabla u_{2} \| \nabla \bar{u}\right| \\
& -2<f\left(u_{1}\right)-f\left(u_{2}\right), \bar{u}>_{L^{p, q}}-2\left(\left((\bar{\eta})^{\prime}, \bar{\eta}\right)\right)_{\mu}  \tag{3.38}\\
\leq & -2 m|\nabla \bar{u}|^{2}+2 m|\nabla \bar{u}|^{2}+\frac{1}{2 m} L_{a}^{2}(R)|l|^{2}|\bar{u}|^{2}\left\|u_{2}\right\|^{2} \\
& -2<f\left(u_{1}\right)-f\left(u_{2}\right), \bar{u}>_{L^{p, q}}-2\left(\left((\bar{\eta})^{\prime}, \bar{\eta}\right)\right)_{\mu} \\
\leq & \frac{1}{2 m} L_{a}^{2}(R)|l|^{2}\|\bar{z}\|_{\mathcal{H}}^{2}\left\|u_{2}\right\|^{2}-2<f\left(u_{1}\right)-f\left(u_{2}\right), \bar{u}>_{L^{p, q}}-2\left(\left((\bar{\eta})^{\prime}, \bar{\eta}\right)\right)_{\mu}
\end{align*}
$$

where $<\cdot, \cdot>_{L^{p, q}}$ is the duality between $L^{2 p}$ and $L^{q}$. The previous calculation is obtained formally taking the product in $\mathcal{H}$ between $\bar{z}$ and the difference of $(3.3)$ with $z_{1}$ and $z_{2}$ in place of $z$, and it can be made rigorous with the use of mollifiers, see [12, Theorem, Section 2]. In fact, integrating by parts and by the fact that $\mu^{\prime}<0$ (see again [12, Section 2]), we have

$$
2\left(\left((\bar{\eta})^{\prime}, \bar{\eta}\right)\right)_{\mu}=-\lim _{s \rightarrow 0} \mu(s)\left|\nabla \bar{\eta}^{t}(s)\right|^{2}-\int_{0}^{\infty} \mu^{\prime}(s)\left|\nabla \bar{\eta}^{t}(s)\right|^{2} d s \geq 0
$$

Hence, the last term of the right hand side of (3.38) can be neglected.
At last, from (2.7) we know that $f(u)$ is increasing for $|u| \geq M$, for some $M>0$. Fix $t \in(\tau, T]$, and let

$$
\Omega_{1}:=\left\{x \in \Omega:\left|u_{1}(t, x)\right| \leq M \text { and }\left|u_{2}(t, x)\right| \leq M\right\},
$$

and

$$
N=2 \sup _{|s| \leq M}\left|f^{\prime}(s)\right| .
$$

Let $x \in \Omega_{1}$, then we have

$$
2\left|f\left(u_{1}(x)\right)-f\left(u_{2}(x)\right)\right| \leq N|\bar{u}(x)| .
$$

Then, by the monotonicity of $f(u)$ for $|u| \geq M$ and the Poincaré inequality, it follows that

$$
\begin{align*}
-2<f\left(u_{1}\right)-f\left(u_{2}\right), \bar{u}>_{L^{p, q}} & \leq-2 \int_{\Omega_{1}}\left(f\left(u_{1}(x)\right)-f\left(u_{2}(x)\right)\right) \bar{u}(x) d x \\
& \leq \int_{\Omega_{1}} N|\bar{u}(x)|^{2} d x  \tag{3.39}\\
& \leq N\|\bar{z}\|_{\mathcal{H}}^{2} .
\end{align*}
$$

(3.38)-(3.39) imply that

$$
\frac{d}{d t}\|\bar{z}\|_{\mathcal{H}}^{2} \leq\left(\frac{1}{2 m} L_{a}^{2}|l|^{2}\left\|u_{2}\right\|^{2}+N\right)\|\bar{z}\|_{\mathcal{H}}^{2}
$$

The uniqueness and continuous dependence on initial data of solution to problem (3.3) follow from the Gronwall inequality. Till now, we finish the proof of the first assertion.

Step 6. (Further regularity) We are going to study further regularity of $(u, \eta)$. To this end, let us first consider the linear operator $\mathcal{I}: L_{V \cap H^{2}(\Omega)}^{2} \rightarrow L_{\mu}^{2}\left(\mathbb{R}^{+} ; D(V)\right)$ defined by

$$
(\mathcal{I} \phi)(s)=\int_{-s}^{0} \phi(r) d r, \quad s \in \mathbb{R}^{+} .
$$

Then, the operator $\mathcal{I}$ defined above is a linear and continuous mapping. In particular, there exists a positive constant $K_{\mu}$, which is the same as in Lemma 3.1, such that, for any $\phi \in L_{V \cap H^{2}(\Omega)}^{2}$, it holds

$$
\|\mathcal{I} \phi\|_{L_{\mu}^{2}\left(\mathbb{R}^{+} ; D(A)\right)}^{2} \leq K_{\mu}\|\phi\|_{L_{V \cap H^{2}(\Omega)}^{2}}^{2} .
$$

Next, multiplying $(2.5)_{1}$ by $-\Delta u$ with respect to the inner product of $H$, the Laplacian of 2.5$)_{2}$ by $\eta$ with respect to the inner product of $L_{\mu}^{2}\left(\mathbb{R}^{+} ; D(A)\right)$, and adding the two terms, we obtain

$$
\begin{equation*}
\frac{d}{d t}\|z\|_{\mathcal{V}}^{2}+2 a(l(u))|\Delta u|^{2}+2\left(\left(\left(\eta^{t},\left(\eta^{t}\right)^{\prime}\right)\right)\right)_{\mu}=2(-f(u)+g, \Delta u) . \tag{3.40}
\end{equation*}
$$

Since $f$ is a polynomial of odd degree, there exists a constant $d_{0}>0$, such that

$$
\begin{equation*}
f^{\prime}(u) \geq-\frac{d_{0}}{2}, \quad \forall u \in \mathbb{R} \tag{3.41}
\end{equation*}
$$

Then, it follows from the above inequality, $(2.7)$, the Green formula and the Young inequality that

$$
\begin{aligned}
2(f(u), \Delta u) & =2 \int_{\Omega} f_{2 p-1} \Delta u d x-2 \int_{\Omega} f^{\prime}(u) \nabla u \cdot \nabla u d x \\
& \leq \frac{2}{m} f_{2 p-1}^{2}|\Omega|+\frac{m}{2}|\Delta u|^{2}+d_{0}|\nabla u|^{2} .
\end{aligned}
$$

Again by the Young inequality, we have

$$
2(g, \Delta u) \leq \frac{m}{2}|\Delta u|^{2}+\frac{2}{m}|g|^{2} .
$$

Together with (2.2), 3.40 becomes

$$
\begin{equation*}
\frac{d}{d t}\|z\|_{\mathcal{V}}^{2}+m|\Delta u|^{2}+2\left(\left(\left(\eta^{t},\left(\eta^{t}\right)^{\prime}\right)\right)\right)_{\mu} \leq \Theta \tag{3.42}
\end{equation*}
$$

where we have used the notation $\Theta=\frac{2}{m} f_{2 p-1}^{2}|\Omega|+d_{0}|\nabla u|^{2}+\frac{2}{m}|g|^{2}$, which belongs to $L^{1}(\tau, T)$. Under the suitable spatial regularity assumptions on $\eta$, integration by parts in time and using ( $h_{1}$ ), we obtain

$$
\left(\left(\left(\eta^{t},\left(\eta^{t}\right)^{\prime}\right)\right)\right)_{\mu}=-\int_{0}^{\infty} \mu^{\prime}(s)\left|\Delta \eta^{t}(s)\right|^{2} d s \geq 0
$$

Therefore, the term $2\left(\left(\left(\eta^{t},\left(\eta^{t}\right)^{\prime}\right)\right)\right)_{\mu}$ in 3.42) can be neglected, we integrate 3.42) between $\tau$ and $t$, where $t \in(\tau, T)$, which leads to

$$
\begin{equation*}
\|z(t)\|_{\mathcal{V}}^{2}+m \int_{\tau}^{t}|\Delta u(s)|^{2} d s \leq\|z(\tau)\|_{\mathcal{V}}^{2}+\int_{\tau}^{t} \Theta(s) d s \tag{3.43}
\end{equation*}
$$

From the above estimation, we conclude that

$$
\begin{aligned}
& u \in L^{\infty}(\tau, T, V) \cap L^{2}(\tau, T ; D(A)) \\
& \eta \in L^{\infty}\left(\tau, T ; L_{\mu}^{2}\left(\mathbb{R}^{+} ; D(A)\right)\right) .
\end{aligned}
$$

Concerning the assertion (ii) of this theorem, the continuity of $u$ follows again using a slightly modified version of [19, Lemma III.1.2]. The continuity of $\eta$ can be proved mimicking the idea of the proof of Step 4 of $(i)$, with $D(A)$ in place of $V$. The proof of this theorem is complete.

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