# ON UNIFORM ATTRACTORS FOR NON-AUTONOMOUS $p$-LAPLACIAN EQUATION WITH A DYNAMIC BOUNDARY CONDITION 

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#### Abstract

In this paper, we consider the non-autonomous p-Laplacian equation with a dynamic boundary condition. The existence and structure of a compact uniform attractor in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ are established for the case of time-dependent internal force $h(t)$. While the nonlinearity $f$ and the boundary nonlinearity $g$ are dissipative for large values without restriction on the growth order of the polynomial.


## 1. Introduction

In this paper, we study the dynamical behavior of solutions of the following non-autonomous parabolic equation with nonlinear dynamic boundary condition:

$$
\begin{cases}u_{t}-\operatorname{div}\left(|\nabla u|^{p-2} \nabla u\right)+f(u)=h(x, t) & \text { in } \Omega,  \tag{1.1}\\ u_{t}+|\nabla u|^{p-2} \partial_{n} u+g(u)=0 & \text { on } \Gamma, \\ u(\tau)=u_{\tau} & \text { in } \bar{\Omega},\end{cases}
$$

where $\Omega$ is a bounded domain of $\mathbb{R}^{n}$ with a smooth boundary $\Gamma$ and $p \geq 2$. $h(x, t) \in L_{\text {loc }}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$. The functions $f$ and $g \in C^{1}(\mathbb{R}, \mathbb{R})$, satisfy the following

[^0]conditions:
\[

$$
\begin{align*}
k_{1}^{\prime}|s|^{q_{1}}-k_{1} & \leq f(s) s \leq k_{2}^{\prime}|s|^{q_{1}}+k_{2}, \quad q_{1} \geq p,  \tag{1.2}\\
k_{3}^{\prime}|s|^{q_{2}}-k_{3} & \leq g(s) s \leq k_{4}^{\prime}|s|^{q_{2}}+k_{4}, \quad q_{2} \geq 2,  \tag{1.3}\\
f^{\prime}(s) & \geq-l \quad \text { and } \quad g^{\prime}(s) \geq-m, \tag{1.4}
\end{align*}
$$
\]

where $l, m>0, k_{i}, k_{i}^{\prime}>0, i=1,2,3,4$.
The dynamic boundary condition arises in hydrodynamics and the heat transfer theory, it is very natural in many mathematical models, such as heat transfer in a solid in contact with a moving fluid, thermoelastic distortion, diffusion phenomena, heat transfer in two medium, problems in fluid dynamics etc. (see [7], [8], [10], [11] and references therein).

Recently, the reaction-diffusion equation with a dynamic boundary condition has been studied by many authors. For example, in [4], [15] considered the phasefield systems with coupling dynamic boundary conditions. Some estimates of convergence rate of the solutions has been obtained in [22], [23].

In this paper, the operator $\Delta_{p}$ in (1.1) denotes the $p$-Laplacian operator. It is obvious that for the case $p=2$, the equation $(1.1)_{1}$ will become the reactiondiffusion equation. For the case of Dirichlet boundary condition, recently, Song et al. [21] obtained the existence of a uniform attractor in $H_{0}^{1}(\Omega)$, where the compactness in $H_{0}^{1}(\Omega)$ was verified by using of the compactness of $L^{q_{1}}(\Omega)$.

As for the reaction-diffusion equation with a dynamic boundary condition, Fan and Zhong [12] obtained the existence of a global attractor in $\left(H^{1}(\Omega) \cap\right.$ $\left.L^{q_{1}}(\Omega)\right) \times L^{q_{1}}(\Gamma)$ under some additional conditions. For the non-autonomous case, in [1], the authors proved the existence of a weak solution, and established the existence of a pullback attractor. [24] proved the existence of a uniform attractor in $L^{q_{1}}(\Omega) \times L^{q_{1}}(\Gamma)$. The authors in [26], [25] considered the long-time behavior of the reaction-diffusion equation with nonlinear boundary condition and competing nonlinearities.

On the other hand, for the $p$-Laplacian equation, Carvalho, Cholewa and Dlotko gave a detailed discussion about Dirichlet boundary condition in [2], and then they proved the existence of $\left(L^{2}(\Omega), L^{2}(\Omega)\right)$-global attractor, see [6]. In Carvalho and Gentile [3], the authors obtained that the corresponding semigroup has a $\left(L^{2}(\Omega), W_{0}^{1, p}(\Omega)\right)$-global attractor.

However, the long time behavior about the $p$-Laplacian equation with dynamic boundary is less discussed, especially for the non-autonomous systems. In this case of autonomous systems, Gal et al. [13], [14] presented the general results about the well-posedness and the asymptotic behavior.

Our main goal of this paper is to study the long-time behavior of solutions of problem (1.1)-(1.4) by the theory of uniform attractors.

The existence and structure of a uniform attractor for the problem (1.1)-(1.4) in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ has been verified.

For the existence of a uniform attractor, as in the autonomous case, some kind of compactness of the family of processes is a key ingredient. In our paper, the growth orders of nonlinear terms $f(u)$ and $g(u)$ have no further restrictions and the solutions have not higher regularities, one can not obtain the compactness of the process in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ by embedding theorem. Furthermore, due to the dynamic boundary conditions, the compactness of the process in $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma)$ apparently can not be obtained, namely, it seems to be difficult to obtain the asymptotic compactness in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ through the compactness of $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma)$ (as that in [21]). Therefore, some new ideas and methods seem to be needed.

In this paper, we testify the uniform asymptotic compactness in $W^{1, p}(\Omega) \times$ $W^{1-1 / p, p}(\Gamma)$ only based on the compactness in $L^{2}(\Omega) \times L^{2}(\Gamma)$ and without any compactness in $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma), q_{1}, q_{2}>2$.

At the same time, we use the closed process to obtain the structure of the uniform attractor, see more details in [24] (see Pata and Zelik [20] for autonomous case).

For convenience, in what follows, we use the notation $\|\cdot\|$ and $\|\cdot\|_{\Gamma}$ stand for the norm in $L^{2}(\Omega)$ and $L^{2}(\Gamma),\langle\cdot, \cdot\rangle$ and $\langle\cdot, \cdot\rangle_{\Gamma}$ stand for the inner product in $L^{2}(\Omega)$ and $L^{2}(\Gamma)$, respectively. $|e|$ denotes the Lebesgue measure of $e$, while $C, C_{i}$ denote general positive constants, $i=1,2, \ldots$, which will be different in different estimates.

Hereafter, we also assume $2 \leq p<N$.
For the case $p \geq N$, the embeddings $W^{1, p}(\Omega) \hookrightarrow L^{s_{1}}(\Omega)$ and $W^{1, p}(\Omega) \hookrightarrow$ $W^{1-1 / p, p}(\Gamma) \hookrightarrow L^{s_{2}}(\Gamma)$ hold for any $s_{1}, s_{2} \in[1, \infty)$, which make the nonlinear terms $f(\cdot)$ and $g(\cdot)$ to be trivial terms.

This paper is organized as follows: in Section 2, we give some preparations for our consideration; in Section 3, the existence and structure of a uniform attractor in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ is obtained.

## 2. Preliminaries

In this section, we first recall some basic concepts about non-autonomous systems, we refer to [5] for more details.

Let $X$ be a Banach space, and $\Sigma$ be a parameter set.
The operators $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ are said to be a family of processes in $X$ with symbol space $\Sigma$ if for any $\sigma \in \Sigma$

$$
\begin{aligned}
U_{\sigma}(t, s) \circ U_{\sigma}(s, \tau) & =U_{\sigma}(t, \tau), & & \text { for all } t \geq s \geq \tau, \tau \in \mathbb{R}, \\
U_{\sigma}(\tau, \tau) & =\mathrm{Id}, & & \text { for all } \tau \in \mathbb{R} .
\end{aligned}
$$

Let $\{T(s)\}_{s \geq 0}$ be the translation semigroup on $\Sigma$, we say that a family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ satisfies the translation identity if

$$
\begin{aligned}
U_{\sigma}(t+s, \tau+s) & =U_{T(s) \sigma}(t, \tau), & & \text { for all } \sigma \in \Sigma, t \geq \tau, \tau \in \mathbb{R}, s \geq 0 \\
T(s) \Sigma & =\Sigma, & & \text { for all } s \geq 0
\end{aligned}
$$

By $\mathcal{B}(X)$ we denote the collection of all bounded sets of $X$ and $\mathbb{R}_{\tau}=\{t \in \mathbb{R}$, $t \geq \tau\}$.

Definition 2.1 ([5]). A bounded set $B_{0} \in \mathcal{B}(X)$ is said to be a bounded uniformly (w.r.t. $\sigma \in \Sigma$ ) absorbing set for $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ if for any $\tau \in \mathbb{R}$ and $B \in \mathcal{B}(X)$ there exists $T_{0}=T_{0}(B, \tau)$ such that $\bigcup_{\sigma \in \Sigma} U_{\sigma}(t ; \tau) B \subset B_{0}$ for all $t \geq T_{0}$.

Definition 2.2 ([5]). A set $\mathcal{A} \subset X$ is said to be uniformly (w.r.t. $\sigma \in \Sigma$ ) attracting for the family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ if for any fixed $\tau \in \mathbb{R}$ and any $B \in \mathcal{B}(X)$

$$
\lim _{t \rightarrow+\infty}\left(\sup _{\sigma \in \Sigma} \operatorname{dist}\left(U_{\sigma}(t ; \tau) B ; \mathcal{A}\right)\right)=0
$$

here dist $(\cdot, \cdot)$ is the usual Hausdorff semidistance in $X$ between two sets.
Definition 2.3 ([5]). A closed set $\mathcal{A}_{\Sigma} \subset X$ is said to be the uniform (w.r.t. $\sigma \in \Sigma$ ) attractor of the family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ if it is uniformly (w.r.t. $\sigma \in \Sigma$ ) attracting (attracting property) and contained in any closed uniformly (w.r.t. $\sigma \in \Sigma$ ) attracting set $\mathcal{A}^{\prime}$ of the family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}: \mathcal{A}_{\Sigma} \subseteq \mathcal{A}^{\prime}$ (minimality property).

Definition 2.4 ([5]). A function $\varphi$ is said to be translation bounded in $L_{\text {loc }}^{2}(\mathbb{R} ; X)$, if

$$
\|\varphi\|_{b}^{2}=\sup _{t \in \mathbb{R}} \int_{t}^{t+1}\|\varphi\|_{X}^{2} d s<+\infty
$$

Denote by $L_{b}^{2}(\mathbb{R} ; X)$ the set of all translation bounded functions in $L_{\text {loc }}^{2}(\mathbb{R} ; X)$.
The next is an estimate of the $p$-Laplacian operator, e.g. see [9] for the proof.
Lemma 2.5. Let $p \geq 2$. Then there exists constant $K>0$ such that for any $a, b \in \mathbb{R}^{n}$,

$$
\left.\left.\langle | a\right|^{p-2} a-|b|^{p-2} b, a-b\right\rangle \geq K|a-b|^{p}
$$

where $K$ depends only on $p$ and $n ;\langle\cdot, \cdot\rangle$ denotes the inner product of $\mathbb{R}^{n}$.

## 3. Uniform attractor in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$

Since $\Omega \subset \mathbb{R}^{n}$ is a bounded domain with smooth boundary $\Gamma$, we define the Sobolev spaces $W^{k, p}(\Omega)$ and $W^{k, p}(\Gamma)$ to be, respectively, the completion of $C^{k}(\bar{\Omega})$ and $C^{k}(\Gamma)$, with respect to the norm

$$
\|u\|_{W^{k, p}(\Omega)}:=\sum_{0 \leq|\alpha| \leq k}\left(\int_{\Omega}\left|\nabla^{\alpha} u\right|^{p} d x\right)^{1 / p}
$$

and

$$
\|u\|_{W^{k, p}(\Gamma)}:=\sum_{j=0}^{k}\left(\int_{\Gamma}\left|\nabla_{\Gamma}^{j} u\right|^{p} d S\right)^{1 / p}
$$

Here, $d x$ denotes the Lebesgue measure on $\Omega$ and $d S$ denotes the natural surface measure on $\Gamma$. For $p \in(1, \infty)$, we define the fractional order Sobolev space
$W^{1-1 / p, p}(\Gamma):=\left\{u \in L^{p}(\Gamma): \int_{\Gamma} \int_{\Gamma}\left(\frac{|u(x)-u(y)|}{|x-y|^{1-1 / p}}\right)^{p} \frac{1}{|x-y|^{N-1}} d S_{x} d S_{y}<\infty\right\}$.
Moreover, since $W^{1, p}(\Omega) \hookrightarrow W^{1-1 / p, p}(\Gamma)$, one has that the norms on $W^{1, p}(\Omega) \times$ $W^{1-1 / p, p}(\Gamma)$ and $W^{1, p}(\Omega)$ are equivalent.

Next, as in [14], we introduce the following rigorous notion of weak solution to our problem.

Definition 3.1 ([1], [14]). The pair of functions $(u(t), v(t))$ is said to be a weak solution of (1.1), if $v(t)=u(t)_{\mid \Gamma}$ in the trace sense, for almost every $t \in(\tau, T)$, for any $\tau, T \in \mathbb{R}, T>\tau$, it satisfying:

$$
\left\{\begin{array}{l}
u(t) \in \mathcal{C}\left([\tau, \infty) ; L^{2}(\Omega)\right) \cap L_{\mathrm{loc}}^{p}\left(\tau, T ; W^{1, p}(\Omega)\right), \\
v(t) \in \mathcal{C}\left([\tau, \infty) ; L^{2}(\Gamma)\right) \cap L_{\mathrm{loc}}^{p}\left(\tau, T ; W^{1-1 / p, p}(\Gamma)\right),
\end{array}\right.
$$

and for all $\sigma \in W^{1, p}(\Omega)$ (hence, $\sigma_{\mid \Gamma} \in W^{1-1 / p, p}(\Gamma)$ ) and for almost every $t \in$ $(\tau, T)$, the following relation holds:

$$
\begin{align*}
\left\langle\partial_{t} u(t), \sigma\right\rangle+\left\langle\partial_{t} v(t), \sigma_{\mid \Gamma}\right\rangle_{\Gamma}+ & \left.\left.\langle | \nabla u\right|^{p-2} \nabla u(t), \nabla \sigma\right\rangle  \tag{3.1}\\
& +\langle f(u(t)), \sigma\rangle+\left\langle g(v(t)), \sigma_{\mid \Gamma}\right\rangle_{\Gamma}=\langle h(t), \sigma\rangle .
\end{align*}
$$

Moreover, in the space $L^{2}(\Omega) \times L^{2}(\Gamma)$, we have $u(\tau)=u_{\tau}, v(\tau)=v_{\tau}$.
Follows the well posedness result in [1], [14], we have the following result and the time-dependent terms make no essential complications.

Theorem 3.2 ([1], [14]). Let $\Omega$ be a bounded domain of $\mathbb{R}^{n}$ with smooth boundary $\Gamma, h(t)$ is translation bounded in $L_{\mathrm{loc}}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right), f$ and $g$ satisfy (1.2)(1.4). Then for any initial data $\left(u_{\tau}, v_{\tau}\right) \in L^{2}(\Omega) \times L^{2}(\Gamma)$, the problem (1.1) has
a unique solution $(u(t), v(t))$. Moreover, $\left(u_{\tau}, v_{\tau}\right) \mapsto(u(t), v(t))$ is continuous on $L^{2}(\Omega) \times L^{2}(\Gamma)$.

We now define the symbol space $\Sigma$ for (1.1). Taking a fixed symbol $\sigma_{0}=$ $h_{0}, h_{0} \in L_{b}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$. We denote by $L_{\text {loc }}^{2, w}\left(\mathbb{R} ; L^{2}(\Omega)\right)$ the space $L_{\text {loc }}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$ endowed with local weak convergence topology.

Set $\Sigma_{0}=\left\{h_{0}(s+h) \mid h \in \mathbb{R}\right\}$, and let $\Sigma$ be the closure of $\Sigma_{0}$ in $L_{\text {loc }}^{2, w}\left(\mathbb{R} ; L^{2}(\Omega)\right)$. Thus, from Theorem 3.2, we know that the problem (1.1)-(1.4) is well posed for all $\sigma(s) \in \Sigma$ and generates a family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ given by the formula:

$$
U_{\sigma}(t, \tau)\left(u_{\tau}, v_{\tau}\right)=(u(t), v(t))
$$

where $(u(t), v(t))$ is the solution of (1.1)-(1.4) and $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ satisfies (2.1)-(2.2). At the same time, due to the unique solvability, we know $\left\{U_{\sigma}(t, \tau)\right.$, $\sigma \in \Sigma\}$ satisfies the translation identity (2.3)-(2.4).

Then, we prove the existence of an uniformly (w.r.t. $\sigma \in \Sigma$ ) bounded absorbing set in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$. The proof is basically same as in [24], and for the sake of completeness, we replicate it here.

Theorem 3.3. Assume that $h(t)$ is translation bounded in $L_{\mathrm{loc}}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$, $f$ and $g$ satisfy (1.2)-(1.3). Then the family of processes $\left\{U_{\sigma}(t, \tau), \sigma \in \Sigma\right\}$ corresponding to (1.1) has a bounded uniformly (w.r.t. $\sigma \in \Sigma$ ) absorbing set $B_{0}$ in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$.

Proof. The following estimates can be deduced by a formal argument, this can be justified by means of the approximation procedure devised in the [14, Theorem 2.6]. Taking $\sigma=u(t)$ and $\sigma_{\mid \Gamma}=v(t)$ in (3.1), we obtain that

$$
\begin{align*}
\frac{1}{2} \frac{d}{d t} \int_{\Omega}|u|^{2} d x & +\frac{1}{2} \frac{d}{d t} \int_{\Gamma}|v|^{2} d S+\frac{1}{2} \int_{\Omega}|\nabla u|^{p} d x  \tag{3.2}\\
& +k_{1}^{\prime} \int_{\Omega}|u|^{p} d x+k_{3}^{\prime} \int_{\Gamma}|v|^{q} d S \leq C+\frac{1}{4 \delta} \int_{\Omega}\left|h_{0}(t, x)\right|^{2} d x
\end{align*}
$$

this implies that

$$
\begin{aligned}
\frac{1}{2} \frac{d}{d t} \int_{\Omega}|u|^{2} d x+\frac{1}{2} \frac{d}{d t} \int_{\Gamma}|v|^{2} d S+\frac{1}{2} \int_{\Omega}|\nabla u|^{p} d x & +C\left(\int_{\Omega}|u|^{2} d x+\int_{\Gamma}|v|^{2} d S\right) \\
& \leq C+\frac{1}{4 \delta} \int_{\Omega}\left|h_{0}(t, x)\right|^{2} d x
\end{aligned}
$$

Using the Gronwall lemma, we know that there exist positive constants $T_{0}>\tau$ and $\alpha>0$, such that

$$
\begin{equation*}
\|u(t)\|^{2}+\|u(t)\|_{\Gamma}^{2} \leq \alpha, \quad \text { for any } t \geq T_{0}, \sigma \in \Sigma \tag{3.3}
\end{equation*}
$$

Then let $F(s)=\int_{0}^{s} f(\tau) d \tau, G(s)=\int_{0}^{s} g(\tau) d \tau$. Using (1.2)-(1.3) again, from (3.2) we deduce that

$$
\begin{aligned}
\frac{d}{d t}\left(\int_{\Omega}|u|^{2} d x+\int_{\Gamma}|v|^{2} d S\right)+\int_{\Omega}|\nabla u|^{p} d x+C_{f}^{\prime} & \int_{\Omega} F(u) d x+C_{g}^{\prime} \int_{\Gamma} G(v) d S \\
& \leq C+\frac{1}{2 \delta} \int_{\Omega}\left|h_{0}(t, x)\right|^{2} d x
\end{aligned}
$$

Integrating the inequality above from $t$ to $t+1$, and combining (3.3), it follows that for any $t \geq T_{0}$, we have

$$
\begin{align*}
\int_{t}^{t+1}\left(\int_{\Omega}|\nabla u|^{p} d x+C_{f}^{\prime} \int_{\Omega} F(u) d x\right. & \left.+C_{g}^{\prime} \int_{\Gamma} G(v) d S\right) d s  \tag{3.4}\\
& \leq C+\frac{1}{2 \delta} \int_{t}^{t+1}\left\|h_{0}(s)\right\|^{2} d s \leq M_{1}
\end{align*}
$$

where the constant $M_{1}$ depends on $|\Omega|, S(\Gamma), \alpha,\|h(t)\|_{b}^{2}$.
On the other hand, taking $\sigma=\partial_{t} u(t)$ and $\sigma_{\mid \Gamma}=\partial_{t} v(t)$ in (3.1), we obtain

$$
\begin{align*}
\int_{\Omega}\left|u_{t}\right|^{2} d x+\int_{\Gamma}\left|v_{t}\right|^{2} d S+\frac{1}{p} \frac{d}{d t} & \|\nabla u\|^{p}+\frac{d}{d t}\left(\int_{\Omega} F(u) d x+\int_{\Gamma} G(v) d S\right)  \tag{3.5}\\
& =\int_{\Omega} h_{0}(t) u_{t} d x \leq \frac{1}{2}\left\|h_{0}(t)\right\|^{2}+\frac{1}{2}\left\|u_{t}\right\|^{2}
\end{align*}
$$

so we obtain

$$
\begin{equation*}
\frac{d}{d t}\left(\|\nabla u\|^{p}+p \int_{\Omega} F(u) d x+p \int_{\Gamma} G(v) d S\right) \leq \frac{p}{2}\left\|h_{0}(t)\right\|^{2} \tag{3.6}
\end{equation*}
$$

Combining (3.4) and (3.6), by the uniformly Gronwall lemma, we have

$$
\begin{equation*}
\|\nabla u\|^{p}+p \int_{\Omega} F(u) d x+p \int_{\Gamma} G(v) d S \leq \rho_{0}, \quad \text { for any } t \geq T_{0}+1, \sigma \in \Sigma \tag{3.7}
\end{equation*}
$$

where $\rho_{0}$ depends on $|\Omega|, S(\Gamma), M_{1},\|h(t)\|_{b}^{2}$. From (3.7), we obtain that for any $t \geq T_{0}+1, \sigma \in \Sigma$, there exists a positive constant $\rho$ depending on $|\Omega|, S(\Gamma), M_{1}$, $\|h(t)\|_{b}^{2}$, such that

$$
\|\nabla u(t)\|^{p}+\|u(t)\|_{L^{q_{1}}(\Omega)}+\|v(t)\|_{L^{q_{2}}(\Gamma)} \leq \rho, \quad \text { for any } t \geq T_{0}+1, \sigma \in \Sigma
$$

As mentioned in [14], since $W^{1, p}(\Omega) \hookrightarrow W^{1-1 / p, p}(\Gamma)$, one has that the norms on $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ and $W^{1, p}(\Omega)$ are equivalent. The proof is complete.

Note that, $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ is compactly embedded into $L^{2}(\Omega) \times L^{2}(\Gamma)$. From Theorem 3.3, the existence of a uniform attractor in $L^{2}(\Omega) \times L^{2}(\Gamma)$ can be obtained immediately.

Corollary 3.4. Under the assumption of Theorem 3.3, the family of processes $\left\{U_{\sigma}(t, \tau)\right\}, \sigma \in \Sigma$ corresponding to (1.1) has a compact uniform (w.r.t. $\sigma \in \Sigma)$ attractor $\mathcal{A}_{\Sigma 0}$ in $L^{2}(\Omega) \times L^{2}(\Gamma)$.

Then, we will give some a priori estimates about $u_{t}$. In what follows, we always denote the weak differential of $h(t)$ with respect to $t$ by $h^{\prime}(t)$.

Lemma 3.5. Let $h(t)$ and $h^{\prime}(t)$ be translation bounded in $L_{\text {loc }}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$, $f$ and $g$ satisfy (1.2)-(1.4), then for any $\tau \in \mathbb{R}$ and any bounded subset $B \subset$ $L^{2}(\Omega) \times L^{2}(\Gamma)$, there exist two positive constants $T=T(B, \tau)>\tau$ and $M_{2}$, such that

$$
\int_{\Omega}\left|u_{t}(s)\right|^{2} d x+\int_{\Gamma}\left|v_{t}(s)\right|^{2} d S \leq M_{2} \quad \text { for all } s \geq T, \quad\left(u_{\tau}, v_{\tau}\right) \in B, \sigma \in \Sigma
$$

where

$$
u_{t}(s)=\left.\frac{d}{d t}\left(U_{\sigma}(t, \tau) u_{\tau}\right)\right|_{t=s} \quad \text { and } \quad v_{t}(s)=\left.\frac{d}{d t}\left(U_{\sigma}(t, \tau) v_{\tau}\right)\right|_{t=s}
$$

$M_{2}$ is a positive constant which depends on $|\Omega|, S(\Gamma), \rho,\|h(t)\|_{b}^{2},\left\|h^{\prime}(t)\right\|_{b}^{2}$.
Proof. Our estimates can be justified by means of the approximation procedure, where we proceed formally. By differentiating (1.1) with external force $h_{0}(t)$ in the time and denoting $\theta=u_{t}, \varrho=v_{t}$, we have

$$
\begin{align*}
\left\langle\partial_{t} \theta, \sigma\right\rangle & \left.+\left\langle\partial_{t} \varrho, \sigma_{\mid \Gamma}\right\rangle_{\Gamma}+\left.\langle | \nabla u\right|^{p-2} \nabla \theta, \nabla \sigma\right\rangle  \tag{3.8}\\
& \left.+\left.(p-2)\langle | \nabla u\right|^{p-4}(\nabla u \cdot \nabla \theta) \nabla u, \nabla \sigma\right\rangle \\
& +\left\langle f^{\prime}(u) \theta, \sigma\right\rangle+\left\langle g^{\prime}(v) \varrho, \sigma_{\mid \Gamma}\right\rangle_{\Gamma}=\langle h(t), \sigma\rangle
\end{align*}
$$

for all $\sigma \in W^{1, p}(\Omega)$ and $\sigma_{\mid \Gamma} \in W^{1-1 / p, p}(\Gamma)$, almost everywhere in $(\tau, \infty)$, where "." denotes the dot product in $\mathbb{R}^{n}, \varrho(t):=\theta(t)_{\mid \Gamma}$.

Taking $\sigma=\theta$ and $\sigma_{\mid \Gamma}=\varrho$ in (3.8), we obtain that

$$
\left.\begin{array}{l}
\frac{1}{2} \frac{d}{d t} \int_{\Omega}|\theta|^{2} d x+\frac{1}{2} \frac{d}{d t} \int_{\Gamma}|\varrho|^{2} d S+\int_{\Omega}|\nabla u|^{p-2}|\nabla \theta|^{2} d x \\
\quad+(p-2) \int_{\Omega}|\nabla u|^{p-4}(\nabla u \cdot \nabla \theta)^{2} d x+\int_{\Omega} f^{\prime}(u) \theta^{2} d x
\end{array}\right)=\int_{\Gamma} g^{\prime}(v) \varrho^{2} d S
$$

From (1.4), this yields

$$
\begin{aligned}
& \frac{1}{2} \frac{d}{d t} \int_{\Omega}|\theta|^{2} d x+\frac{1}{2} \frac{d}{d t} \int_{\Gamma}|\varrho|^{2} d S \\
& \quad+\int_{\Omega}|\nabla u|^{p-2}|\nabla \theta|^{2} d x+(p-2) \int_{\Omega}|\nabla u|^{p-4}(\nabla u \cdot \nabla \theta)^{2} d x \\
& \quad \leq l \int_{\Omega}|\theta|^{2} d x+m \int_{\Gamma}|\varrho|^{2} d S+\frac{1}{2} \int_{\Omega}|\theta|^{2} d x+\frac{1}{2}\left\|h_{0}^{\prime}(t)\right\|^{2},
\end{aligned}
$$

so we have

$$
\begin{equation*}
\frac{d}{d t} \int_{\Omega}|\theta|^{2} d x+\frac{d}{d t} \int_{\Gamma}|\varrho|^{2} d S \leq C\left(\int_{\Omega}|\theta|^{2} d x+\int_{\Gamma}|\varrho|^{2} d S\right)+\left\|h_{0}^{\prime}(t)\right\|^{2} \tag{3.9}
\end{equation*}
$$

On the other hand, integrating (3.5) from $t$ to $t+1$, and using (3.7), we have

$$
\begin{equation*}
\int_{t}^{t+1}\left(\int_{\Omega}|\theta|^{2} d x+\int_{\Gamma}|\varrho|^{2} d S\right) \leq \widetilde{C} \tag{3.10}
\end{equation*}
$$

where $\widetilde{C}$ depends on $|\Omega|, S(\Gamma), M,\|h(t)\|_{b}^{2}$. Combining (3.9)-(3.10), and using the uniform Gronwall lemma, we get

$$
\int_{\Omega}\left|u_{t}(s)\right|^{2} d x+\int_{\Gamma}\left|v_{t}(s)\right|^{2} d S \leq M_{2} \quad \text { for all } s \geq T, \quad\left(u_{\tau}, v_{\tau}\right) \in B, \sigma \in \Sigma
$$

where $M_{2}$ depends on $|\Omega|, S(\Gamma), M,\|h(t)\|_{b}^{2},\left\|h^{\prime}(t)\right\|_{b}^{2}$.
Finally, the following theorem gives the existence and structure of an uniform attractor in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ :

Theorem 3.6. Assume that $h(t) \in L^{\infty}\left(\mathbb{R} ; L^{2}(\Omega)\right)$ and $h^{\prime}(t)$ is translation bounded in $L_{\mathrm{loc}}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right), f$ and $g$ satisfy (1.2)-(1.4). Then the family of processes $\left\{U_{\sigma}(t, \tau)\right\}, \sigma \in \Sigma$ corresponding to (1.1) has a compact uniform (w.r.t. $\sigma \in \Sigma)$ attractor $\mathcal{A}_{\Sigma 1}$ in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ and $\mathcal{A}_{\Sigma 1}$ satisfies:

$$
\mathcal{A}_{\Sigma 1}=\omega_{0, \Sigma}\left(B_{0}\right)=\bigcup_{\sigma \in \Sigma} \mathcal{K}_{\sigma}(s), \quad \text { for all } s \in \mathbb{R}
$$

where $\mathcal{K}_{\sigma}(s)$ is the section at $t=s$ of the kernel $\mathcal{K}_{\sigma}$ of the process $\left\{U_{\sigma}(t, \tau)\right\}$ with symbol $\sigma$.

Proof. Let $B_{0}$ be a $\left(W^{1, p}(\Omega) \cap L^{q_{1}}(\Omega) \times W^{1-1 / p, p}(\Gamma) \cap L^{q_{2}}(\Gamma)\right)$-bounded uniformly (w.r.t. $\sigma \in \Sigma$ ) absorbing set obtained in Theorem 3.3, then we need only to show that:
(3.11) for any $\left\{\left(u_{\tau_{n}}, v_{\tau_{n}}\right)\right\} \subset B_{0},\left\{\sigma_{n}\right\} \subset \Sigma$ and $t_{n} \rightarrow \infty$,

$$
\left\{\left(U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) u_{\tau_{n}}, U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) v_{\tau_{n}}\right)\right\}_{n=1}^{\infty}
$$

is precompact in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$.
Thanks to Corollary 3.4, we know that $\left\{\left(U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) u_{\tau_{n}}, U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) v_{\tau_{n}}\right)\right\}_{n=1}^{\infty}$ is precompact in $L^{2}(\Omega) \times L^{2}(\Gamma)$. Without loss of generality, we assume that $\left\{\left(U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) u_{\tau_{n}}, U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) v_{\tau_{n}}\right)\right\}_{n=1}^{\infty}$ is a Cauchy sequence in $L^{2}(\Omega) \times L^{2}(\Gamma)$.

Next, we prove that $\left\{\left(U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) u_{\tau_{n}}, U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) v_{\tau_{n}}\right)\right\}_{n=1}^{\infty}$ is a Cauchy sequence in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$.

Denote by $u_{n}^{\sigma_{n}}\left(t_{n}\right):=U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) u_{\tau_{n}}, v_{n}^{\sigma_{n}}\left(t_{n}\right):=U_{\sigma_{n}}\left(t_{n}, \tau_{n}\right) v_{\tau_{n}}$, from Lemma 2.5 , which is the property of $p$-Laplacian operator when $p \geq 2$, and using (1.4) again, we know that there exists a constant $c>0$, such that

$$
\begin{aligned}
c\left(\| u_{n}^{\sigma_{n}}\left(t_{n}\right)\right. & \left.-u_{m}^{\sigma_{m}}\left(t_{m}\right)\left\|_{W^{1, p}(\Omega)}^{p}+\right\| v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right) \|_{W^{1-1 / p, p}(\Gamma)}^{p}\right) \\
\leq & \int_{\Omega}\left|\frac{d}{d t} u_{n}^{\sigma_{n}}\left(t_{n}\right)-\frac{d}{d t} u_{m}^{\sigma_{m}}\left(t_{m}\right)\right|\left|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right| \\
& +\int_{\Gamma}\left|\frac{d}{d t} v_{n}^{\sigma_{n}}\left(t_{n}\right)-\frac{d}{d t} v_{m}^{\sigma_{m}}\left(t_{m}\right)\right|\left|v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right)\right| \\
& +\int_{\Omega}\left|\sigma_{n}-\sigma_{m}\right|\left|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right| \\
& +l\left\|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|^{2}+m\left\|v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|_{\Gamma}^{2},
\end{aligned}
$$

which implies that

$$
\begin{aligned}
c\left(\| u_{n}^{\sigma_{n}}\left(t_{n}\right)\right. & \left.-u_{m}^{\sigma_{m}}\left(t_{m}\right)\left\|_{W^{1, p}(\Omega)}^{p}+\right\| v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right) \|_{W^{1-1 / p, p}(\Gamma)}^{p}\right) \\
\leq & \left\|\frac{d}{d t} u_{n}^{\sigma_{n}}\left(t_{n}\right)-\frac{d}{d t} u_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|\left\|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right\| \\
& +\left\|\frac{d}{d t} v_{n}^{\sigma_{n}}\left(t_{n}\right)-\frac{d}{d t} v_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|_{\Gamma}\left\|v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|_{\Gamma} \\
& +\left\|\sigma_{n}-\sigma_{m}\right\|\left\|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right\| \\
& +l\left\|u_{n}^{\sigma_{n}}\left(t_{n}\right)-u_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|^{2}+m\left\|v_{n}^{\sigma_{n}}\left(t_{n}\right)-v_{m}^{\sigma_{m}}\left(t_{m}\right)\right\|_{\Gamma}^{2},
\end{aligned}
$$

which, combining with Theorem 3.3 and Lemma 3.5, and since the norms on $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ and $W^{1, p}(\Omega)$ are equivalent, we have (3.11) immediately. Then, we use the closed process to obtain the structure of $\mathcal{A}_{\Sigma 1}$ in $W^{1, p}(\Omega) \times$ $W^{1-1 / p, p}(\Gamma)$, see more details in [24] (see Pata and Zelik [20] for autonomous case).

Remark 3.7. Note that, the growth orders of nonlinear terms $f(u)$ and $g(u)$ have no further restrictions and the solutions have not higher regularities, one can not obtain the compactness in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ by an embedding theorem. Furthermore, it seems difficult to obtain the compactness in $W^{1, p}(\Omega) \times$ $W^{1-1 / p, p}(\Gamma)$ through the compactness of $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma)$ (as that in [12], [21]).

REMARK 3.8. In this paper, the compactness in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ was verified only by using of the compactness in $L^{2}(\Omega) \times L^{2}(\Gamma)$ and without any compactness in $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma), q_{1}, q_{2}>2$. This implies that the compactness of the process in $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$ did not depend on the compactness of the process in $L^{q_{1}}(\Omega) \times L^{q_{2}}(\Gamma), q_{1}, q_{2}>2$, i.e. did not depend on the growth orders of nonlinear terms $f$ and $g$ only if the nonlinear terms $f$ and $g$ satisfy a very weak condition that $f^{\prime} \geq-l, g^{\prime} \geq-m$.

Remark 3.9. Using the argument of the closed process (see more details in Pata and Zelik [20]), we can easily obtain the structure of the uniform attractors.

Remark 3.10. In Theorem 3.6, the assumption $h(x, t) \in L^{\infty}\left(\mathbb{R} ; L^{2}(\Omega)\right)$ is only needed to guarantee the uniform asymptotic compactness in $W^{1, p}(\Omega) \times$ $W^{1-1 / p, p}(\Gamma)$. In fact, if we are only concerned with the existence of the uniform attractor in $L^{q_{1}}(\Omega) \times L^{q_{1}}(\Gamma)$, then we only assume that $h(x, t) \in L_{n}^{2}\left(\mathbb{R} ; L^{2}(\Omega)\right)$ (i.e. normal, see [24] for more details).

Remark 3.11. As for the autonomous case of (1.1), that is $h(x, t)=h(x)$, under the assumption that $h(x) \in L^{2}(\Omega)$, the method in Section 3 also is valid, and the main result - Theorem 3.6 also holds.

REMARK 3.12. In this paper, we study the asymptotic behavior of the solutions of problem (1.1) by the concept of uniform attractors. For the nonautonomous dynamical systems, the theory of pullback attractors is also a good tool to describe the long time behavior of the solutions, see more detail in [16], etc. When considered for pullback attractors, the external forces $h(x, t)$ usually only satisfy some weaker condition than $h(x, t)$ of this paper (see, e.g. [19]), and it seems difficult to directly apply the method of this paper for obtaining the $W^{1, p}(\Omega) \times W^{1-1 / p, p}(\Gamma)$-compactness, especially, we can not perform as that in Lemma 3.5 to derive the estimates of $u_{t}$ and $v_{t}$.

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