

**DIMENSION OF ATTRACTORS  
AND INVARIANT SETS  
IN REACTION DIFFUSION EQUATIONS**

MARTINO PRIZZI

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ABSTRACT. Under fairly general assumptions, we prove that every compact invariant set  $\mathcal{I}$  of the semiflow generated by the semilinear reaction diffusion equation

$$\begin{aligned}u_t + \beta(x)u - \Delta u &= f(x, u), & (t, x) \in [0, +\infty[ \times \Omega, \\u &= 0, & (t, x) \in [0, +\infty[ \times \partial\Omega\end{aligned}$$

in  $H_0^1(\Omega)$  has finite Hausdorff dimension. Here  $\Omega$  is an arbitrary, possibly unbounded, domain in  $\mathbb{R}^3$  and  $f(x, u)$  is a nonlinearity of subcritical growth. The nonlinearity  $f(x, u)$  needs not to satisfy any dissipativeness assumption and the invariant subset  $\mathcal{I}$  needs not to be an attractor. If  $\Omega$  is regular,  $f(x, u)$  is dissipative and  $\mathcal{I}$  is the global attractor, we give an explicit bound on the Hausdorff dimension of  $\mathcal{I}$  in terms of the structure parameter of the equation.

### 1. Introduction

In this paper we consider the reaction diffusion equation

$$(1.1) \quad \begin{aligned}u_t + \beta(x)u - \Delta u &= f(x, u), & (t, x) \in [0, +\infty[ \times \Omega, \\u &= 0, & (t, x) \in [0, +\infty[ \times \partial\Omega.\end{aligned}$$

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Here  $\Omega$  is an arbitrary (possibly unbounded) open set in  $\mathbb{R}^3$ ,  $\beta(x)$  is a potential such that the operator  $-\Delta + \beta(x)$  is positive, and  $f(x, u)$  is a nonlinearity of subcritical growth (i.e. of polynomial growth strictly less than five).

The assumptions on  $\beta(x)$  and  $f(x, u)$  will be made more precise in Section 2 below. Under such assumptions, equation (1.1) generates a local semiflow  $\pi$  in the space  $H_0^1(\Omega)$ . Suppose that the semiflow  $\pi$  admits a compact invariant set  $\mathcal{I}$  (i.e.  $\pi(t, \mathcal{I}) = \mathcal{I}$  for all  $t \geq 0$ ). We do not make any structure assumption on the nonlinearity  $f(x, u)$  and therefore we do not assume that  $\mathcal{I}$  is the global attractor of equation (1.1): for example,  $\mathcal{I}$  can be an unstable invariant set detected by Conley index arguments (see e.g. [16]).

Our aim is to prove that  $\mathcal{I}$  has finite Hausdorff dimension and to give an explicit estimate of its dimension. The first results concerning the dimension of invariant sets of dynamical systems are due to Mallet–Paret [14] and Mañé [15]. For a comprehensive study of the subject, see e.g. [6], [12], [20], [23].

When  $\Omega$  is a bounded domain and  $f(x, u)$  satisfies suitable dissipativeness conditions, the existence of a finite dimensional compact global attractor for (1.1) is a classical achievement (see e.g. [6], [12], [23]). When  $\Omega$  is unbounded, new difficulties arise due to the lack of compactness of the Sobolev embeddings. These difficulties can be overcome in several ways: by introducing weighted spaces (see e.g. [5], [9]), by developing suitable tail-estimates (see e.g. [24], [17]), by exploiting comparison arguments (see e.g. [3]).

Concerning the finite dimensionality of the attractor, in [5], [9], [24] and other similar works the potential  $\beta(x)$  is always assumed to be just a positive constant. In [4] Arrieta et al. considered for the first time the case of a sign-changing potential. In their results the invariant set  $\mathcal{I}$  does not need to be an attractor; however they need to make some structure assumptions on  $f(x, u)$  which essentially resemble the conditions ensuring the existence of the global attractor. Moreover, in [4] the invariant set is a-priori assumed to be bounded in the  $L^\infty$ -norm. In concrete situations, such a-priori estimate can be obtained through elliptic regularity combined with some comparison argument. This in turn requires to make some regularity assumption on the boundary of  $\Omega$ .

In this paper we do not make any structure assumption on the nonlinearity  $f(x, u)$ , neither do we assume  $\partial\Omega$  to be regular. Our only assumption is that the mapping  $h \mapsto (\partial_u f(x, 0))_+ h$  has to be a relatively form compact perturbation of  $-\Delta + \beta(x)$ . This can be achieved, e.g. by assuming that  $\partial_u f(x, 0)$  can be estimated from above by some positive  $L^r$  function,  $r > 3/2$ . Under this assumption, we shall prove that  $\mathcal{I}$  has finite Hausdorff dimension. Also, we give an explicit estimate of the dimension of  $\mathcal{I}$ , involving the number  $\mathcal{N}$  of negative eigenvalues of the operator  $-\Delta + \beta(x) - \partial_u f(x, 0)$ . When  $\Omega$  has a regular boundary, we can explicitly estimate  $\mathcal{N}$  by mean of Cwikel–Lieb–Rozenblum inequality (see [21]);

as a consequence, if we also assume that  $f(x, u)$  is dissipative, we recover the result of Arrieta et al. [4].

The paper is organized as follows. In Section 2 we introduce notations, we state the main assumptions and we collect some preliminaries about the semiflow generated by equation (1.1). In Section 3 we prove that the semiflow generated by equation (1.1) is uniformly  $L^2$ -differentiable on any compact invariant set  $\mathcal{I}$ . In Section 4 we recall the definition of Hausdorff dimension and we prove that any compact invariant set  $\mathcal{I}$  has finite Hausdorff dimension in  $L^2(\Omega)$  as well as in  $H_0^1(\Omega)$ . In Section 5 we compute the number of negative eigenvalues of the operator  $-\Delta + \beta(x) - \partial_u f(x, 0)$  by mean of Cwickel–Lieb–Rozenblum inequality. In Section 6 we specialize our result to the case of a dissipative equation and we recover the result of Arrieta et al. [4].

The results contained in this paper continue to hold if one replaces  $-\Delta$  with the general second order elliptic operator in divergence form

$$-\sum_{i,j=1}^3 \partial_{x_i}(a_{ij}(x)\partial_{x_j}).$$

### 2. Notation, preliminaries and remarks

Let  $\sigma \geq 1$ . We denote by  $L_u^\sigma(\mathbb{R}^N)$  the set of measurable functions  $\omega: \mathbb{R}^N \rightarrow \mathbb{R}$  such that

$$|\omega|_{L_u^\sigma} := \sup_{y \in \mathbb{R}^N} \left( \int_{B(y)} |\omega(x)|^\sigma dx \right)^{1/\sigma} < \infty,$$

where, for  $y \in \mathbb{R}^N$ ,  $B(y)$  is the open unit cube in  $\mathbb{R}^N$  centered at  $y$ .

In this paper we assume throughout that  $N = 3$ , and we fix an open (possibly unbounded) set  $\Omega \subset \mathbb{R}^3$ . We denote by  $M_B$  the constant of the Sobolev embedding  $H^1(B) \subset L^6(B)$ , where  $B$  is any open unit cube in  $\mathbb{R}^3$ . Moreover, for  $2 \leq q \leq 6$ , we denote by  $M_q$  the constant of the Sobolev embedding  $H^1(\mathbb{R}^3) \subset L^q(\mathbb{R}^3)$ .

**PROPOSITION 2.1.** *Let  $\sigma > 3/2$  and let  $\omega \in L_u^\sigma(\mathbb{R}^3)$ . Set  $\rho := 3/2\sigma$ . Then, for every  $\varepsilon > 0$  and for every  $u \in H_0^1(\Omega)$ ,*

$$\int_{\Omega} |\omega(x)||u(x)|^2 dx \leq |\omega|_{L_u^\sigma}(\rho\varepsilon M_B^2 |u|_{H^1}^2 + (1 - \rho)\varepsilon^{-\rho/(1-\rho)} |u|_{L^2}^2).$$

Moreover, for every  $u \in H_0^1(\Omega)$ ,

$$\int_{\Omega} |\omega(x)||u(x)|^2 dx \leq M_B^{2\rho} |\omega|_{L_u^\sigma} |u|_{H^1}^{2\rho} |u|_{L^2}^{2(1-\rho)}.$$

**PROOF.** See the proof of Lemma 3.3 in [18]. □

Let  $\beta \in L^{\sigma}_{\mathbb{u}}(\mathbb{R}^3)$ , with  $\sigma > 3/2$ . Let us consider the following bilinear form defined on the space  $H_0^1(\Omega)$ :

$$a(u, v) := \int_{\Omega} \nabla u(x) \cdot \nabla v(x) \, dx + \int_{\Omega} \beta(x) u(x) v(x) \, dx, \quad u, v \in H_0^1(\Omega).$$

Our first assumption is the following:

**HYPOTHESIS 2.2.** *There exists  $\lambda_1 > 0$  such that*

$$(2.1) \quad \int_{\Omega} |\nabla u(x)|^2 \, dx + \int_{\Omega} \beta(x) |u(x)|^2 \, dx \geq \lambda_1 |u|_{L^2}^2, \quad u \in H_0^1(\Omega).$$

**REMARK 2.3.** Conditions on  $\beta(x)$  under which Hypothesis 2.2 is satisfied are expounded e.g. in [1], [2].

As a consequence of (2.1) and Proposition 2.1, we have:

**PROPOSITION 2.4.** *There exist two positive constants  $\lambda_0$  and  $\Lambda_0$  such that*

$$\lambda_0 |u|_{H^1}^2 \leq \int_{\Omega} |\nabla u(x)|^2 \, dx + \int_{\Omega} \beta(x) |u(x)|^2 \, dx \leq \Lambda_0 |u|_{H^1}^2, \quad u \in H_0^1(\Omega).$$

The constants  $\lambda_0$  and  $\Lambda_0$  can be computed explicitly in terms of  $\lambda_1$ ,  $M_B$  and  $|\beta|_{L^{\sigma}_{\mathbb{u}}}$ .

**PROOF.** Cf. Lemma 4.2 in [17].  $\square$

It follows from Proposition 2.4 that the bilinear form  $a(\cdot, \cdot)$  defines a scalar product in  $H_0^1(\Omega)$ , equivalent to the standard one. According to the results of Section 4 in [17],  $a(\cdot, \cdot)$  induces a positive selfadjoint operator  $A$  in the space  $L^2(\Omega)$ .  $A$  is uniquely determined by the relation

$$\langle Au, v \rangle_{L^2} = a(u, v), \quad u \in D(A), \quad v \in H_0^1(\Omega).$$

Notice that  $Au = -\Delta u + \beta u$  in the sense of distributions, and  $u \in D(A)$  if and only if  $-\Delta u + \beta u \in L^2(\Omega)$ . Set  $X := L^2(\Omega)$ , and let  $(X^{\alpha})_{\alpha \in \mathbb{R}}$  be the scale of fractional power spaces associated with  $A$  (see Section 2 in [17] for a short, self-contained, description of this scale of spaces). Here we just recall that  $X^0 = L^2(\Omega)$ ,  $X^1 = D(A)$ ,  $X^{1/2} = H_0^1(\Omega)$  and  $X^{-\alpha}$  is the dual of  $X^{\alpha}$  for  $\alpha \in ]0, +\infty[$ . For  $\alpha \in ]0, +\infty[$ , the space  $X^{\alpha}$  is a Hilbert space with respect to the scalar product

$$\langle u, v \rangle_{X^{\alpha}} := \langle A^{\alpha} u, A^{\alpha} v \rangle_{L^2}, \quad u, v \in X^{\alpha}.$$

Also, the space  $X^{-\alpha}$  is a Hilbert space with respect to the scalar product  $\langle \cdot, \cdot \rangle_{X^{-\alpha}}$  dual to the scalar product  $\langle \cdot, \cdot \rangle_{X^{\alpha}}$ , i.e.

$$\langle u', v' \rangle_{X^{-\alpha}} = \langle R_{\alpha}^{-1} u', R_{\alpha}^{-1} v' \rangle_{X^{\alpha}}, \quad u, v \in X^{-\alpha},$$

where  $R_{\alpha}: X^{\alpha} \rightarrow X^{-\alpha}$  is the Riesz isomorphism  $u \mapsto \langle \cdot, u \rangle_{X^{\alpha}}$ . Finally, for every  $\alpha \in \mathbb{R}$ ,  $A$  induces a selfadjoint operator  $A_{(\alpha)}: X^{\alpha+1} \rightarrow X^{\alpha}$ , such that  $A_{(\alpha')}$  is

an extension of  $A_{(\alpha)}$  whenever  $\alpha' \leq \alpha$ , and  $D(A_{(\alpha)}^\beta) = X^{\alpha+\beta}$  for  $\beta \in [0, 1]$ . If  $\alpha \in [0, 1/2]$ ,  $u \in X^{1-\alpha}$  and  $v \in X^{1/2} \subset X^\alpha$ , then

$$\langle v, A_{(-\alpha)}u \rangle_{(X^\alpha, X^{-\alpha})} = \langle u, v \rangle_{X^{1/2}} = a(u, v).$$

LEMMA 2.5. *Let  $(X^\alpha)_{\alpha \in \mathbb{R}}$  be as above.*

- (a) *If  $p \in [2, 6[$ , then  $X^\alpha \subset L^p(\Omega)$  for  $\alpha \in ]3(p-2)/4p, 1/2]$ . Accordingly, if  $q \in ]6/5, 2]$ , then  $L^q(\Omega) \subset X^{-\alpha}$  for  $\alpha \in ]3(2-q)/4q, 1/2]$ .*
- (b) *If  $\sigma > 3/2$  and  $\omega \in L_u^\sigma(\Omega)$ , then the assignment  $u \mapsto \omega u$  defines a bounded linear map from  $X^{1/2}$  to  $X^{-\alpha}$  for  $\alpha \in ]3/4\sigma, 1/2]$ .*

PROOF. See Lemmas 5.1 and 5.2 and the proof of Proposition 5.3 in [17].  $\square$

Our second assumption is the following:

HYPOTHESIS 2.6.

- (a)  *$f: \Omega \times \mathbb{R} \rightarrow \mathbb{R}$  is such that, for every  $u \in \mathbb{R}$ ,  $f(\cdot, u)$  is measurable and, for almost every  $x \in \Omega$ ,  $f(x, \cdot)$  is of class  $C^2$ ;*
- (b)  *$f(\cdot, 0) \in L^q(\Omega)$ , with  $6/5 < q \leq 2$  and  $\partial_u f(\cdot, 0) \in L_u^\sigma(\mathbb{R}^3)$ , with  $\sigma > 3/2$ ;*
- (c) *there exist constants  $C$  and  $\gamma$ , with  $C > 0$  and  $2 \leq \gamma < 3$  such that  $|\partial_{uu} f(x, u)| \leq C(1 + |u|^\gamma)$ . Notice that, in view of Young's inequality, the requirement  $\gamma \geq 2$  is not restrictive.*

We introduce the Nemitski operator  $\hat{f}$  which associates with every function  $u: \Omega \rightarrow \mathbb{R}$  the function  $\hat{f}(u)(x) := f(x, u(x))$ .

PROPOSITION 2.7. *Assume  $f$  satisfies Hypothesis (2.6). Let  $\alpha$  be such that*

$$\frac{1}{2} > \alpha > \max \left\{ \frac{\gamma - 1}{4}, \frac{3}{4} \frac{2 - q}{q}, \frac{3}{4\sigma} \right\}.$$

*Then the assignment  $u \mapsto \mathbf{f}(u)$ , where*

$$\langle v, \mathbf{f}(u) \rangle_{(X^\alpha, X^{-\alpha})} := \int_{\Omega} \hat{f}(u)(x)v(x) dx,$$

*defines a map  $\mathbf{f}: X^{1/2} \rightarrow X^{-\alpha}$  which is Lipschitzian on bounded sets.*

PROOF. See the proof of Proposition 5.3 in [17].  $\square$

Setting  $\mathbf{X} := X^{-\alpha}$  and  $\mathbf{A} := A_{(-\alpha)}$ , we have that  $\mathbf{X}^{\alpha+1/2} = X^{1/2}$ . We can rewrite equation (1.1) as an abstract parabolic problem in the space  $\mathbf{X}$ , namely

$$(2.2) \quad \dot{u} + \mathbf{A}u = \mathbf{f}(u).$$

By results in [11], equation (2.2) has a unique *mild solution* for every initial datum  $u_0 \in \mathbf{X}^{\alpha+1/2} = H_0^1(\Omega)$ , satisfying the *variation of constants* formula

$$u(t) = e^{-\mathbf{A}t}u_0 + \int_0^t e^{-\mathbf{A}(t-s)}\mathbf{f}(u(s)) ds, \quad t \geq 0.$$

It follows that (2.2) generates a local semiflow  $\pi$  in the space  $H_0^1(\Omega)$ . Moreover, if  $u(\cdot): ]0, T[ \rightarrow \mathbf{X}^{\alpha+1/2}$  is a mild solution of (2.2), then  $u(t)$  is differentiable into  $\mathbf{X}^{\alpha+1/2} = H_0^1(\Omega)$  for  $t \in ]0, T[$ , and it satisfies equation (2.2) in  $\mathbf{X} = X^{-\alpha} \subset H^{-1}(\Omega)$ . In particular,  $u(\cdot)$  is a *weak solution* of (1.1).

Assume now that  $\mathcal{I} \subset H_0^1(\Omega)$  is a compact invariant set for the semiflow  $\pi$  generated by (2.2). If  $\mathcal{B}$  is a Banach space such that  $H_0^1(\Omega) \subset \mathcal{B}$ , we define

$$|\mathcal{I}|_{\mathcal{B}} := \max\{|u|_{\mathcal{B}} \mid u \in \mathcal{I}\}.$$

We end this section with a technical lemma that will be used later.

LEMMA 2.8. *For every  $T > 0$  there exists a constant  $L(T)$  such that, whenever  $u_0$  and  $v_0 \in \mathcal{I}$ , setting  $u(t) := \pi(t, u_0)$  and  $v(t) := \pi(t, v_0)$ ,  $t \geq 0$ , the following estimate holds:*

$$|u(t) - v(t)|_{H^1} \leq L(T)t^{-(\alpha+1/2)}|u_0 - v_0|_{L^2}, \quad t \in ]0, T].$$

The constant  $L(T)$  depends only on  $|\mathcal{I}|_{H^1}$ , and on the constants of Hypotheses 2.2 and 2.6.

PROOF. We have

$$u(t) - v(t) = e^{-\mathbf{A}t}(u_0 - v_0) + \int_0^t e^{-\mathbf{A}(t-s)}(\mathbf{f}(u(s)) - \mathbf{f}(v(s))) ds;$$

it follows that

$$\begin{aligned} & |u(t) - v(t)|_{\mathbf{X}^{\alpha+1/2}} \\ & \leq t^{-(\alpha+1/2)}|u_0 - v_0|_{\mathbf{X}} + \int_0^t (t-s)^{-(\alpha+1/2)}|\mathbf{f}(u(s)) - \mathbf{f}(v(s))|_{\mathbf{X}} ds \\ & \leq t^{-(\alpha+1/2)}|u_0 - v_0|_{\mathbf{X}} + \int_0^t (t-s)^{-(\alpha+1/2)}C(|\mathcal{I}|_{H^1})|u(s) - v(s)|_{\mathbf{X}^{\alpha+1/2}} ds. \end{aligned}$$

By Henry's inequality [11, Theorem 7.1.1], this implies that

$$|u(t) - v(t)|_{\mathbf{X}^{\alpha+1/2}} \leq L(T)t^{-(\alpha+1/2)}|u_0 - v_0|_{\mathbf{X}}, \quad t \in ]0, T],$$

and the thesis follows.  $\square$

### 3. Uniform differentiability

In this section we prove some technical results which will allow us to apply the methods of [23] for proving finite dimensionality of compact invariant sets. We assume throughout that  $\mathcal{I} \subset H_0^1(\Omega)$  is a compact invariant set of the semiflow  $\pi$  generated by equation (2.2).

LEMMA 3.1. *There exists a constant  $K$  such that, whenever  $u_0$  and  $v_0 \in \mathcal{I}$ , setting  $u(t) := \pi(t, u_0)$  and  $v(t) := \pi(t, v_0)$ ,  $t \geq 0$ , the following estimate holds:*

$$|u(t) - v(t)|_{L^2}^2 + \lambda_0 \int_0^t |u(s) - v(s)|_{H^1}^2 ds \leq e^{Kt} |u_0 - v_0|_{L^2}^2.$$

The constant  $K$  depends only on  $|\mathcal{I}|_{H^1}$ , on  $\lambda_0$  and  $\Lambda_0$  (see Proposition 2.4), on  $|\partial_u f(\cdot, 0)|_{L_u^\sigma}$ , and on the constants  $C$  and  $\gamma$  (see Hypothesis 2.6).

PROOF. Set  $z(t) = u(t) - v(t)$ . Then

$$\begin{aligned} \frac{1}{2} \frac{d}{dt} |z(t)|_{L^2}^2 + \int_{\Omega} |\nabla z(t)(x)|^2 dx + \int_{\Omega} \beta(x) |z(t)(x)|^2 dx \\ = \int_{\Omega} (f(x, u(t)(x)) - f(x, v(t)(x))) z(t)(x) dx. \end{aligned}$$

It follows from Proposition 2.4 and Hypothesis 2.6 that

$$\begin{aligned} \frac{1}{2} \frac{d}{dt} |z(t)|_{L^2}^2 + \lambda_0 |z(t)|_{H^1}^2 &\leq \int_{\Omega} |\partial_u f(x, 0)| |z(t)(x)|^2 dx \\ &\quad + C' \int_{\Omega} (1 + |u(t)(x)|^{\gamma+1} + |v(t)(x)|^{\gamma+1}) |z(t)(x)|^2 dx \\ &\leq \int_{\Omega} |\partial_u f(x, 0)| |z(t)(x)|^2 dx + C' |z(t)|_{L^2}^2 \\ &\quad + C' (|u(t)|_{L^6}^{\gamma+1} + |v(t)|_{L^6}^{\gamma+1}) |z(t)|_{L^{12/(5-\gamma)}}^2, \end{aligned}$$

where  $C'$  is a constant depending only on  $C$  and  $\gamma$ . Notice that  $2 < 12/(5-\gamma) < 6$ . Therefore, by interpolation, we get that for every  $\varepsilon > 0$  there exists a constant  $c_\varepsilon > 0$  such that

$$(3.1) \quad |z(t)|_{L^{12/(5-\gamma)}}^2 \leq \varepsilon |z(t)|_{H^1}^2 + c_\varepsilon |z(t)|_{L^2}^2.$$

Now (3.1) and Proposition 2.1 imply that, for every  $\varepsilon > 0$ , there exists a constant  $C'_\varepsilon$ , depending on  $C'$ ,  $|\mathcal{I}|_{H^1}$  and  $\varepsilon$ , such that

$$(3.2) \quad \frac{1}{2} \frac{d}{dt} |z(t)|_{L^2}^2 + \lambda_0 |z(t)|_{H^1}^2 \leq \varepsilon |z(t)|_{H^1}^2 + C'_\varepsilon |z(t)|_{L^2}^2.$$

Now choosing  $\varepsilon = \lambda_0/2$  and multiplying (3.2) by  $e^{-2C'_\varepsilon t}$  we get

$$(3.3) \quad \frac{d}{dt} (e^{-2C'_\varepsilon t} |z(t)|_{L^2}^2) + \lambda_0 e^{-2C'_\varepsilon t} |z(t)|_{H^1}^2 \leq 0.$$

Integrating (3.3) we obtain the thesis. □

Let  $\bar{u}(\cdot): \mathbb{R} \rightarrow H_0^1(\Omega)$  be a full bounded solution of (2.2) such that  $\bar{u}(t) \in \mathcal{I}$  for  $t \in \mathbb{R}$ . Let us consider the non autonomous linear equation

$$(3.4) \quad \begin{aligned} u_t + \beta(x)u - \Delta u &= \partial_u f(x, \bar{u}(t))u, & (t, x) \in [0, +\infty[ \times \Omega, \\ u &= 0, & (t, x) \in [0, +\infty[ \times \partial\Omega. \end{aligned}$$

We introduce the following bilinear form defined on on the space  $H_0^1(\Omega)$ :

$$(3.5) \quad \begin{aligned} a(t; u, v) &:= \int_{\Omega} \nabla u(x) \cdot \nabla v(x) dx + \int_{\Omega} \beta(x)u(x)v(x) dx \\ &\quad - \int_{\Omega} \partial_u f(x, \bar{u}(t)(x))u(x)v(x) dx, \quad u, v \in H_0^1(\Omega). \end{aligned}$$

PROPOSITION 3.2. *There exist constants  $\kappa_i > 0$ ,  $i = 1, \dots, 4$ , such that:*

- (a)  $|a(t; u, v)| \leq \kappa_1 |u|_{H^1} |v|_{H^1}$ ,  $u, v \in H_0^1(\Omega)$ ,  $t \in \mathbb{R}$ ;
- (b)  $|a(t; u, u)| \geq \kappa_2 |u|_{H^1}^2 - \kappa_3 |u|_{L^2}^2$ ,  $u \in H_0^1(\Omega)$ ,  $t \in \mathbb{R}$ ;
- (c)  $|a(t; u, v) - a(s; u, v)| \leq \kappa_4 |t - s| |u|_{H^1} |v|_{H^1}$ ,  $u, v \in H_0^1(\Omega)$ ,  $t, s \in \mathbb{R}$ .

PROOF. Properties (a) and (b) follow from Hypothesis 2.6 and Proposition 2.1. In order to prove point (c), we first observe that, by Theorem 3.5.2 in [11] (and its proof), the function  $\bar{u}(\cdot)$  is differentiable into  $H_0^1(\Omega)$ , with  $|\dot{\bar{u}}(\cdot)|_{H^1} \leq L$ , where  $L$  is a constant depending on  $|\mathcal{I}|_{H^1}$  and on the constants in Hypotheses 2.2 and 2.6. Therefore we have:

$$\begin{aligned} |a(t; u, v) - a(s; u, v)| &\leq \int_{\Omega} |\partial_u f(x, \bar{u}(t)) - \partial_u f(x, \bar{u}(s))| |u(x)| |v(x)| dx \\ &\leq \int_{\Omega} C(1 + |\bar{u}(t)(x)|^\gamma + |\bar{u}(s)(x)|^\gamma) |\bar{u}(t)(x) - \bar{u}(s)(x)| |u(x)| |v(x)| dx \\ &\leq C'(1 + |\bar{u}(t)|_{H^1}^\gamma + |\bar{u}(s)|_{H^1}^\gamma) |\bar{u}(t) - \bar{u}(s)|_{H^1} |u|_{H^1} |v|_{H^1} \\ &\leq C'(1 + 2|\mathcal{I}|_{H^1}^\gamma) L |t - s| |u|_{H^1} |v|_{H^1}, \end{aligned}$$

and the proof is complete.  $\square$

Now let  $A(t)$  be the self-adjoint operator determined by the relation

$$(3.6) \quad \langle A(t)u, v \rangle_{L^2} = a(t; u, v), \quad u \in D(A(t)), \quad v \in H_0^1(\Omega).$$

We can apply Theorem 3.1 in [10] and get:

PROPOSITION 3.3. *There exists a two parameter family of bounded linear operators  $U(t, s): L^2(\Omega) \rightarrow L^2(\Omega)$ ,  $t \geq s$ , such that:*

- (a)  $U(s, s) = I$  for all  $s \in \mathbb{R}$ , and  $U(t, s)U(s, r) = U(t, r)$  for all  $t \geq s \geq r$ ;
- (b)  $U(t, s)h_0 \in D(A(t))$  for all  $h_0 \in L^2(\Omega)$  and  $t > s$ ;
- (c) for every  $h_0 \in L^2(\Omega)$  and  $s \in \mathbb{R}$ , the map  $t \mapsto U(t, s)h_0$  is differentiable into  $L^2(\Omega)$  for  $t > s$ , and

$$\frac{\partial}{\partial t} U(t, s)h_0 = -A(t)U(t, s)h_0.$$

In particular,  $U(t, s)h_0$  is a weak solution of (3.4).

Given  $\bar{u}_0 \in \mathcal{I}$ , we take a full bounded solution  $\bar{u}(\cdot)$  of (2.2), whose trajectory is contained in  $\mathcal{I}$ , and such that  $\bar{u}(0) = \bar{u}_0$ . Then we define

$$(3.7) \quad \mathcal{U}(\bar{u}_0; t) := U(t, 0), \quad t \geq 0,$$

where  $U(t, s)$  is the family of operators given by Proposition 3.3. Notice that  $\mathcal{U}(\bar{u}_0; t)$  does not depend on the choice of  $\bar{u}(\cdot)$ , due to forward uniqueness for equation (2.2).

PROPOSITION 3.4. For every  $t \geq 0$ ,

$$\sup_{\bar{u}_0 \in \mathcal{I}} \|\mathcal{U}(\bar{u}_0; t)\|_{\mathcal{L}(L^2, L^2)} < +\infty.$$

PROOF. Let  $\bar{u}_0 \in \mathcal{I}$  and  $h_0 \in L^2(\Omega)$ . Set  $h(t) := \mathcal{U}(\bar{u}_0; t)h_0$ . Then, by property (c) of Proposition 3.3, for  $t > 0$  we have

$$\begin{aligned} \frac{d}{dt} \frac{1}{2} |h(t)|_{L^2}^2 + \int_{\Omega} |\nabla h(t)(x)|^2 dx + \int_{\Omega} \beta(x) |h(t)(x)|^2 dx \\ = \int_{\Omega} \partial_u f(x, \bar{u}(t)(x)) |h(t)(x)|^2 dx, \end{aligned}$$

where  $\bar{u}(\cdot)$  is a full bounded solution of (2.2), whose trajectory is contained in  $\mathcal{I}$ , and such that  $\bar{u}(0) = \bar{u}_0$ . It follows from Hypothesis 2.6 and Propositions 2.1 and 2.4 that for all  $\varepsilon > 0$

$$\begin{aligned} & \frac{d}{dt} \frac{1}{2} |h(t)|_{L^2}^2 + \lambda_0 |h(t)|_{H^1}^2 \\ & \leq \int_{\Omega} \partial_u f(x, 0) |h(t)(x)|^2 dx + \int_{\Omega} (\partial_u f(x, \bar{u}(t)(x)) - \partial_u f(x, 0)) |h(t)(x)|^2 dx \\ & \leq \varepsilon |h(t)|_{H^1}^2 + c_\varepsilon |h(t)|_{L^2}^2 + \int_{\Omega} C(1 + |\bar{u}(t)(x)|^\gamma) |\bar{u}(t)(x)| |h(t)(x)|^2 dx \\ & \leq \varepsilon |h(t)|_{H^1}^2 + c_\varepsilon |h(t)|_{L^2}^2 + \int_{\Omega} C'(1 + |\bar{u}(t)(x)|^{\gamma+1}) |h(t)(x)|^2 dx \\ & \leq \varepsilon |h(t)|_{H^1}^2 + (c_\varepsilon + C') |h(t)|_{L^2}^2 + C' |\bar{u}(t)|_{L^6}^{\gamma+1} |h(t)|_{L^{12/(5-\gamma)}}^2. \end{aligned}$$

Since  $2 < 12/(5 - \gamma) < 6$ , by interpolation we get that for every  $\varepsilon > 0$  there exists a constant  $c'_\varepsilon > 0$  such that

$$|h(t)|_{L^{12/(5-\gamma)}}^2 \leq \varepsilon |h(t)|_{H^1}^2 + c'_\varepsilon |h(t)|_{L^2}^2.$$

Therefore we have

$$(3.8) \quad \frac{d}{dt} \frac{1}{2} |h(t)|_{L^2}^2 + \lambda_0 |h(t)|_{H^1}^2 \leq \varepsilon |h(t)|_{H^1}^2 + C''(\varepsilon, |\mathcal{I}|_{H^1}) |h(t)|_{L^2}^2.$$

Choosing  $\varepsilon = \lambda_0$  and integrating (3.8) we obtain

$$|h(t)|_{L^2}^2 \leq e^{2C''(\lambda_0, |\mathcal{I}|_{H^1})t} |h_0|_{L^2}^2 \quad \square$$

PROPOSITION 3.5. *For every  $t \geq 0$ ,*

$$\lim_{\varepsilon \rightarrow 0} \sup_{\substack{\bar{u}_0, \bar{v}_0 \in \mathcal{I} \\ 0 < |\bar{u}_0 - \bar{v}_0|_{L^2} < \varepsilon}} \frac{|\pi(t, \bar{v}_0) - \pi(t, \bar{u}_0) - \mathcal{U}(\bar{u}_0; t)(\bar{v}_0 - \bar{u}_0)|_{L^2}}{|\bar{v}_0 - \bar{u}_0|_{L^2}} = 0.$$

PROOF. Let  $\bar{u}_0, \bar{v}_0 \in \mathcal{I}$ . Set  $\bar{u}(t) := \pi(t, \bar{u}_0)$ ,  $\bar{v}(t) := \pi(t, \bar{v}_0)$  and  $\theta(t) := \bar{v}(t) - \bar{u}(t) - \mathcal{U}(\bar{u}_0; t)(\bar{v}_0 - \bar{u}_0)$ ,  $t \geq 0$ . A computation using property (c) of Proposition 3.3 shows that, for  $t > 0$ ,

$$\begin{aligned} & \frac{d}{dt} \frac{1}{2} |\theta(t)|_{L^2}^2 + \int_{\Omega} |\nabla \theta(t)(x)|^2 dx + \int_{\Omega} \beta(x) |\theta(t)(x)|^2 dx \\ &= \int_{\Omega} \partial_u f(x, \bar{u}(t)(x)) |\theta(t)(x)|^2 dx \\ &+ \int_{\Omega} (f(x, \bar{v}(t)(x)) - f(x, \bar{u}(t)(x)) - \partial_u f(x, \bar{u}(t)(x))(\bar{v}(t)(x) - \bar{u}(t)(x))) \theta(t)(x) dx. \end{aligned}$$

Therefore, by Proposition 2.4

$$\frac{d}{dt} \frac{1}{2} |\theta(t)|_{L^2}^2 + \lambda_0 |\theta(t)|_{H^1} \leq I_1(t) + I_2(t) + I_3(t),$$

where

$$\begin{aligned} I_1(t) &:= \int_{\Omega} \partial_u f(x, 0) |\theta(t)(x)|^2 dx, \\ I_2(t) &:= \int_{\Omega} (\partial_u f(x, \bar{u}(t)(x)) - \partial_u f(x, 0)) |\theta(t)(x)|^2 dx, \\ I_3(t) &= \int_{\Omega} (f(x, \bar{v}(t)) - f(x, \bar{u}(t)) - \partial_u f(x, \bar{u}(t))(\bar{v}(t) - \bar{u}(t))) \theta(t) dx. \end{aligned}$$

Repeating the same computations of the proof of Proposition 3.4, for  $\varepsilon > 0$  we get

$$I_1(t) + I_2(t) \leq \varepsilon |\theta(t)|_{H^1}^2 + C_1(\varepsilon, |\mathcal{I}|_{H^1}) |\theta(t)|_{L^2}^2.$$

Concerning  $I_3(t)$ , for  $\varepsilon > 0$  we have

$$\begin{aligned} I_3(t) &\leq \int_{\Omega} C(1 + |\bar{u}(t)(x)|^\gamma + |\bar{v}(t)(x)|^\gamma) |\bar{v}(t)(x) - \bar{u}(t)(x)|^2 \theta(t)(x) dx \\ &\leq C |\theta(t)|_{L^6} |\bar{v}(t) - \bar{u}(t)|_{L^{12/5}}^2 + C |\theta(t)|_{L^6} (|\bar{u}(t)|_{L^6}^\gamma \\ &\quad + |\bar{v}(t)|_{L^6}^\gamma) |\bar{v}(t) - \bar{u}(t)|_{L^{12/(5-\gamma)}}^2 \\ &\leq \varepsilon |\theta(t)|_{H^1}^2 + C_2(\varepsilon, |\mathcal{I}|_{H^1}) (|\bar{v}(t) - \bar{u}(t)|_{L^{12/5}}^4 + |\bar{v}(t) - \bar{u}(t)|_{L^{12/(5-\gamma)}}^4) \\ &\leq \varepsilon |\theta(t)|_{H^1}^2 + C_3(\varepsilon, |\mathcal{I}|_{H^1}) (|\bar{v}(t) - \bar{u}(t)|_{H^1} |\bar{v}(t) - \bar{u}(t)|_{L^2}^3 \\ &\quad + |\bar{v}(t) - \bar{u}(t)|_{H^1}^{1+\gamma} |\bar{v}(t) - \bar{u}(t)|_{L^2}^{3-\gamma}) \end{aligned}$$

Choosing  $\varepsilon = \lambda_0/2$ , we get

$$\begin{aligned} & \frac{d}{dt} \frac{1}{2} |\theta(t)|_{L^2}^2 - C_1(\varepsilon, |\mathcal{I}|_{H^1}) |\theta(t)|_{L^2}^2 \\ & \leq C_3(\varepsilon, |\mathcal{I}|_{H^1}) (|\bar{v}(t) - \bar{u}(t)|_{H^1} |\bar{v}(t) - \bar{u}(t)|_{L^2}^3 + |\bar{v}(t) - \bar{u}(t)|_{H^1}^{1+\gamma} |\bar{v}(t) - \bar{u}(t)|_{L^2}^{3-\gamma}) \\ & \leq C_4(\varepsilon, |\mathcal{I}|_{H^1}) (|\bar{v}(t) - \bar{u}(t)|_{L^2}^3 + |\bar{v}(t) - \bar{u}(t)|_{H^1}^2 |\bar{v}(t) - \bar{u}(t)|_{L^2}^{3-\gamma}). \end{aligned}$$

By Lemma 3.1, we get

$$\begin{aligned} & \frac{d}{dt} \frac{1}{2} |\theta(t)|_{L^2}^2 - C_1(\varepsilon, |\mathcal{I}|_{H^1}) |\theta(t)|_{L^2}^2 \\ & \leq C_4(\varepsilon, |\mathcal{I}|_{H^1}) (e^{3Kt} |\bar{v}_0 - \bar{u}_0|_{L^2}^3 + e^{(3-\gamma)Kt} |\bar{v}(t) - \bar{u}(t)|_{H^1}^2 |\bar{v}_0 - \bar{u}_0|_{L^2}^{3-\gamma}). \end{aligned}$$

Writing  $C_1$  for  $C_1(\varepsilon, |\mathcal{I}|_{H^1})$  and  $C_4$  for  $C_4(\varepsilon, |\mathcal{I}|_{H^1})$ , we have

$$\begin{aligned} (3.9) \quad & \frac{d}{dt} \frac{1}{2} (e^{-C_1 t} |\theta(t)|_{L^2}^2) \\ & \leq C_4 (e^{(3K-C_1)t} |\bar{v}_0 - \bar{u}_0|_{L^2}^3 + e^{((3-\gamma)K-C_1)t} |\bar{v}(t) - \bar{u}(t)|_{H^1}^2 |\bar{v}_0 - \bar{u}_0|_{L^2}^{3-\gamma}). \end{aligned}$$

Finally, integrating (3.9), recalling that  $\theta(0) = 0$  and taking into account Lemma 3.1, we get the existence of two increasing functions  $\Phi_1(t)$  and  $\Phi_2(t)$  such that

$$|\theta(t)|_{L^2}^2 \leq \Phi_1(t) |\bar{v}_0 - \bar{u}_0|_{L^2}^3 + \Phi_2(t) |\bar{v}_0 - \bar{u}_0|_{L^2}^{5-\gamma},$$

and the thesis follows. □

#### 4. Dimension of invariant sets

Let  $\mathcal{X}$  be a complete metric space and let  $\mathcal{K} \subset \mathcal{X}$  be a compact set. For  $d \in \mathbb{R}^+$  and  $\varepsilon > 0$  one defines

$$\mu_H(\mathcal{K}, d, \varepsilon) := \inf \left\{ \sum_{i \in I} r_i^d \mid \mathcal{K} \subset \bigcup_{i \in I} B(x_i, r_i), r_i \leq \varepsilon \right\},$$

where the infimum is taken over all the finite coverings of  $\mathcal{K}$  with balls of radius  $r_i \leq \varepsilon$ . Observe that  $\mu_H(\mathcal{K}, d, \varepsilon)$  is a non increasing function of  $\varepsilon$  and  $d$ . The  $d$ -dimensional Hausdorff measure of  $\mathcal{K}$  is by definition

$$\mu_H(\mathcal{K}, d) := \lim_{\varepsilon \rightarrow 0} \mu_H(\mathcal{K}, d, \varepsilon) = \sup_{\varepsilon > 0} \mu_H(\mathcal{K}, d, \varepsilon).$$

One has:

- (1)  $\mu_H(\mathcal{K}, d) \in [0, +\infty]$ ;
- (2) if  $\mu_H(\mathcal{K}, \bar{d}) < \infty$ , then  $\mu_H(\mathcal{K}, d) = 0$  for all  $d > \bar{d}$ ;
- (3) if  $\mu_H(\mathcal{K}, \bar{d}) > 0$ , then  $\mu_H(\mathcal{K}, d) = +\infty$  for all  $d < \bar{d}$ .

The Hausdorff dimension of  $\mathcal{K}$  is the smallest  $d$  for which  $\mu_H(\mathcal{K}, d)$  is finite, i.e.

$$\dim_H(\mathcal{K}) := \inf \{ d > 0 \mid \mu_H(\mathcal{K}, d) = 0 \}.$$

As pointed up in [22], the Hausdorff dimension is in fact an intrinsic metric property of the set  $\mathcal{K}$ . Moreover, if  $\mathcal{Y}$  is another complete metric space and  $\ell: \mathcal{K} \rightarrow \mathcal{Y}$  is a Lipschitzian map, then  $\dim_H(\ell(\mathcal{K})) \leq \dim_H(\mathcal{K})$ .

There is a well developed technique to estimate the Hausdorff dimension of an invariant set of a map or a semigroup. We refer the reader e.g. to [23] and [12]. The geometric idea consists in tracking the evolution of a  $d$ -dimensional volume under the action of the linearization of the semigroup along solutions lying in the invariant set. One looks then for the smallest  $d$  for which any  $d$ -dimensional volume contracts asymptotically as  $t \rightarrow \infty$ .

Let  $\bar{u}_0 \in \mathcal{I}$  and let  $\bar{u}(\cdot): \mathbb{R} \rightarrow H_0^1(\Omega)$  be a full bounded solution of (2.2) such that  $\bar{u}(0) = \bar{u}_0$  and  $\bar{u}(t) \in \mathcal{I}$  for  $t \in \mathbb{R}$ . For  $t \geq 0$ , we denote by  $a_{\bar{u}_0}(t; u, v)$  the bilinear form defined by (3.5), and by  $A_{\bar{u}_0}(t)$  the self-adjoint operator determined by the relation (3.6). Given a  $d$ -dimensional subspace  $E_d$  of  $L^2(\Omega)$ , with  $E_d \subset H_0^1(\Omega)$ , we define the operator  $A_{\bar{u}_0}(t | E_d): E_d \rightarrow E_d$  by

$$\langle A_{\bar{u}_0}(t | E_d)\phi, \psi \rangle_{L^2} := a_{\bar{u}_0}(t; \phi, \psi), \quad \phi, \psi \in E_d.$$

Notice that, if  $E_d \subset D(A_{\bar{u}_0}(t))$ , then one has  $A_{\bar{u}_0}(t | E_d) = P_{E_d} A_{\bar{u}_0}(t) P_{E_d}|_{E_d}$ , where  $P_{E_d}: L^2(\Omega) \rightarrow E_d$  is the  $L^2$ -orthogonal projection onto  $E_d$ . We define

$$\text{Tr}_d(A_{\bar{u}_0}(t)) := \inf_{\substack{E_d \subset H_0^1(\Omega) \\ \dim E_d = d}} \text{Tr}(A_{\bar{u}_0}(t | E_d)).$$

Let  $\bar{u}_0 \in \mathcal{I}$ , let  $d \in \mathbb{N}$  and let  $v_{0,i} \in L^2(\Omega)$ ,  $i = 1, \dots, d$ . Set  $v_i(t) := \mathcal{U}(\bar{u}_0; t)v_{0,i}$ ,  $t \geq 0$ , where  $\mathcal{U}(\bar{u}_0; t)$  is defined by (3.7). We denote by  $G(t)$  the  $d$ -dimensional volume delimited by  $v_1(t), \dots, v_d(t)$  in  $L^2(\Omega)$ , that is

$$G(t) := |v_1(t) \wedge v_2(t) \wedge \dots \wedge v_d(t)|_{\wedge^d L^2} = (\det(\langle v_i(t), v_j(t) \rangle_{L^2})_{ij})^{1/2}.$$

An easy computation using Leibnitz rule and Proposition 3.3 shows that, for  $t > 0$ ,  $G(t)$  satisfies the ordinary differential equation

$$G'(t) = -\text{Tr}(A_{\bar{u}_0}(t | E_d(t))G(t),$$

where  $E_d(t) := \text{span}(v_1(t), \dots, v_d(t))$ . It follows from Propositions 3.4 and 3.5 and from the results in [23, Chapter V] that the Hausdorff dimension  $\dim_H(\mathcal{I})$  of  $\mathcal{I}$  in  $L^2(\Omega)$  is finite and less than or equal to  $d$ , provided

$$\limsup_{t \rightarrow \infty} \sup_{\bar{u}_0 \in \mathcal{I}} \frac{1}{t} \int_0^t -\text{Tr}_d(A_{\bar{u}_0}(s)) ds < 0.$$

In order to prove that  $\dim_H(\mathcal{I}) \leq d$ , we are lead to estimate  $-\text{Tr}_d(A_{\bar{u}_0}(t))$ . To this end, we notice that, whenever  $E_d$  is a  $d$ -dimensional subspace of  $L^2(\Omega)$ , and  $B: E_d \rightarrow E_d$  is a selfadjoint operator, then

$$\text{Tr}(B) = \sum_{i=1}^d \langle B\phi_i, \phi_i \rangle_{L^2},$$

where  $\phi_1, \dots, \phi_d$  is any  $L^2$ -orthonormal basis of  $E_d$ . So let  $E_d \subset H_0^1(\Omega)$  be a  $d$ -dimensional space and let  $\phi_1, \dots, \phi_d$  be an  $L^2$ -orthonormal basis of  $E_d$ . Fix  $0 < \delta < 1$ . It follows that

$$\begin{aligned} & \text{Tr}(A_{\bar{u}_0}(t \mid E_d)) \\ &= \sum_{i=1}^d \left( (1 - \delta) \left( \int_{\Omega} |\nabla \phi_i|^2 dx + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) - \int_{\Omega} \partial_u f(x, 0) |\phi_i|^2 dx \right) \\ & \quad + \delta \sum_{i=1}^d \left( \int_{\Omega} |\nabla \phi_i|^2 dx + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) \\ & \quad + \sum_{i=1}^d \int_{\Omega} (\partial_u f(x, \bar{u}(t)) - \partial_u f(x, 0)) |\phi_i|^2 dx. \end{aligned}$$

We introduce the following bilinear form defined on the space  $H_0^1(\Omega)$ :

$$\begin{aligned} a_{\delta}(u, v) := (1 - \delta) \left( \int_{\Omega} \nabla u(x) \cdot \nabla v(x) dx + \int_{\Omega} \beta(x) u(x) v(x) dx \right) \\ - \int_{\Omega} \partial_u f(x, 0) u(x) v(x) dx, \quad u, v \in H_0^1(\Omega). \end{aligned}$$

Let  $A_{\delta}$  be the self-adjoint operator determined by the relation

$$\langle A_{\delta} u, v \rangle_{L^2} = a_{\delta}(u, v), \quad u \in D(A_{\delta}), v \in H_0^1(\Omega).$$

Given a  $d$ -dimensional subspace  $E_d$  of  $L^2(\Omega)$ , with  $E_d \subset H_0^1(\Omega)$ , we define the operator  $A_{\delta}(E_d): E_d \rightarrow E_d$  by

$$\langle A_{\delta}(E_d) \phi, \psi \rangle_{L^2} := a_{\delta}(\phi, \psi), \quad \phi, \psi \in E_d.$$

It follows that

$$\begin{aligned} \text{Tr}(A_{\bar{u}_0}(t \mid E_d)) &= \text{Tr}(A_{\delta}(E_d)) + \delta \sum_{i=1}^d \left( \int_{\Omega} |\nabla \phi_i|^2 dx \right. \\ & \quad \left. + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) + \sum_{i=1}^d \int_{\Omega} (\partial_u f(x, \bar{u}(t)) - \partial_u f(x, 0)) |\phi_i|^2 dx. \end{aligned}$$

We introduce the *proper values* of the operator  $A_{\delta}$ :

$$\mu_j(A_{\delta}) := \sup_{\psi_1, \dots, \psi_{j-1} \in H_0^1(\Omega)} \inf_{\substack{\psi \in [\psi_1, \dots, \psi_{j-1}]^{\perp} \\ \|\psi\|_{L^2} = 1, \psi \in H_0^1(\Omega)}} a_{\delta}(\psi, \psi) \quad j = 1, 2, \dots$$

We recall (see e.g. Theorem XIII.1 in [19]) that:

PROPOSITION 4.1. *For each fixed  $n$ , either*

- (a) *there are at least  $n$  eigenvalues (counting multiplicity) below the bottom of the essential spectrum of  $A_{\delta}$  and  $\mu_n(A_{\delta})$  is the  $n$ th eigenvalue (counting multiplicity);*

or

- (b)  $\mu_n(A_\delta)$  is the bottom of the essential spectrum and in that case  $\mu_{n+j}(A_\delta) = \mu_n(A_\delta)$ ,  $j = 1, 2, \dots$  and there are at most  $n-1$  eigenvalues (counting multiplicity) below  $\mu_n(A_\delta)$ .  $\square$

Let  $\mu_j(A_\delta(E_d))$ ,  $j = 1, \dots, d$ , be the eigenvalues of  $A_\delta(E_d)$ . By Theorem XIII.3 in [19], we have that

$$\mu_j(A_\delta(E_d)) \geq \mu_j(A_\delta), \quad j = 1, \dots, d.$$

It follows that

$$\begin{aligned} \text{Tr}(A_{\bar{u}_0}(t | E_d)) &\geq \sum_{i=1}^d \mu_i(A_\delta) + \delta \sum_{i=1}^d \left( \int_{\Omega} |\nabla \phi_i|^2 dx + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) \\ &\quad + \sum_{i=1}^d \int_{\Omega} (\partial_u f(x, \bar{u}(t)) - \partial_u f(x, 0)) |\phi_i|^2 dx. \end{aligned}$$

To proceed further, we need to recall the Lieb–Thirring inequality (see [13]).

PROPOSITION 4.2. *Let  $N \in \mathbb{N}$  and let  $p \in \mathbb{R}$ , with  $\max\{N/2, 1\} \leq p \leq 1 + N/2$ . There exists a constant  $K_{p,N} > 0$  such that, if  $\phi_1, \dots, \phi_d \in H^1(\mathbb{R}^N)$  are pairwise  $L^2$ -orthonormal, then*

$$(4.1) \quad \sum_{i=1}^d \int_{\mathbb{R}^N} |\nabla \phi_i(x)|^2 dx \geq \frac{1}{K_{p,N}} \left( \int_{\mathbb{R}^N} \rho(x)^{p/(p-1)} dx \right)^{2(p-1)/N},$$

where  $\rho(x) := \sum_{i=1}^d |\phi_i(x)|^2$ .

Now we have:

LEMMA 4.3. *Let  $\bar{u} \in \mathcal{I}$  and let  $\phi_1, \dots, \phi_d \in H^1(\mathbb{R}^N)$  be pairwise  $L^2$ -orthonormal. Then*

$$\begin{aligned} \delta \sum_{i=1}^d \left( \int_{\Omega} |\nabla \phi_i|^2 dx + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) \\ + \sum_{i=1}^d \int_{\Omega} (\partial_u f(x, \bar{u}(x)) - \partial_u f(x, 0)) |\phi_i|^2 dx \geq -D(\gamma, \lambda_0, \delta, |\mathcal{I}|_{H^1}), \end{aligned}$$

where

$$\begin{aligned} D(\gamma, \lambda_0, \delta, |\mathcal{I}|_{H^1}) &= \frac{5}{2} \left( \frac{3}{5} \frac{2}{\delta \lambda_0} \right)^{3/2} (C|\mathcal{I}|_{L^{5/2}} K_{5/2,3})^{5/2} \\ &\quad + \frac{3-\gamma}{4} \left( \frac{\gamma+1}{4} \frac{2}{\delta \lambda_0} \right)^{(\gamma+1)/(3-\gamma)} (C|\mathcal{I}|_{L^6}^{\gamma+1} K_{6/(\gamma+1),3})^{4/(3-\gamma)}. \end{aligned}$$

PROOF. We observe first that

$$\delta \sum_{i=1}^d \left( \int_{\Omega} |\nabla \phi_i|^2 dx + \int_{\Omega} \beta(x) |\phi_i|^2 dx \right) \geq \delta \lambda_0 \sum_{i=1}^d \int_{\Omega} |\nabla \phi_i|^2 dx.$$

On the other hand,

$$\begin{aligned} & \left| \int_{\Omega} \left( \partial_u f(x, \bar{u}(x)) - \partial_u f(x, 0) \right) \rho(x) dx \right| \\ & \leq \int_{\Omega} C(1 + |\bar{u}|^\gamma) |\bar{u}| |\rho| dx \leq C |\bar{u}|_{L^{5/2}} |\rho|_{L^{5/3}} + C |\bar{u}|_{L^6}^{\gamma+1} |\rho|_{L^{6/(5-\gamma)}}. \end{aligned}$$

By Lieb–Thirring inequality (4.1), we have

$$\begin{aligned} & \left| \int_{\Omega} (\partial_u f(x, \bar{u}(x)) - \partial_u f(x, 0)) \rho(x) dx \right| \\ & \leq C |\mathcal{I}|_{L^{5/2}} K_{5/2,3} \left( \sum_{i=1}^d \int_{\mathbb{R}^N} |\nabla \phi_i|^2 dx \right)^{3/5} \\ & \quad + C |\mathcal{I}|_{L^6}^{\gamma+1} K_{6/(\gamma+1),3} \left( \sum_{i=1}^d \int_{\mathbb{R}^N} |\nabla \phi_i|^2 dx \right)^{(\gamma+1)/4}. \end{aligned}$$

The conclusion follows by a simple application of Young’s inequality. □

Thanks to Lemma 4.3, we finally get:

$$\text{Tr}(A_{\bar{u}_0}(t | E_d)) \geq \sum_{i=1}^d \mu_i(A_\delta) - D(\gamma, \lambda_0, \delta, |\mathcal{I}|_{H^1}).$$

Therefore, in order to conclude that  $\dim_H(\mathcal{I})$  is finite, we are lead to make some assumption which guarantees that  $\sum_{i=1}^d \mu_i(A_\delta)$  can be made positive and as large as we want, by choosing  $d$  sufficiently large. This is equivalent to the fact that the bottom of the essential spectrum of  $A_\delta$  be strictly positive. We make the following assumption:

**HYPOTHESIS 4.4.** *For every  $\varepsilon > 0$  there exists  $V_\varepsilon \in L^r(\Omega)$ ,  $r > 3/2$ ,  $V_\varepsilon \geq 0$ , such that  $\partial_u f(x, 0) \leq V_\varepsilon(x) + \varepsilon$ , for  $x \in \Omega$ .*

We need the following lemmas:

**LEMMA 4.5.** *Let  $r > 3/2$  and let  $V \in L^r(\Omega)$ . If  $r > 3$  let  $p := 2$ ; if  $r \leq 3$  let  $p := 6/5$ . Then the assignment  $u \mapsto Vu$  defines a compact map from  $H_0^1(\Omega)$  to  $L^p(\Omega)$ , and hence to  $H^{-1}(\Omega)$ .*

PROOF. Let  $B \subset H_0^1(\Omega)$  be bounded. If  $\mathcal{B}$  is a Banach space such that  $H_0^1(\Omega) \subset \mathcal{B}$ , we define  $|B|_{\mathcal{B}} := \sup\{|u|_{\mathcal{B}} \mid u \in B\}$ . If  $u \in H_0^1(\Omega)$  we denote by  $\tilde{u}$  its trivial extension to the whole  $\mathbb{R}^3$ . Similarly, we denote by  $\tilde{V}$  the trivial

extension of  $V$  to  $\mathbb{R}^3$ . For  $k > 0$ , let  $\chi_k$  be the characteristic function of the set  $\{x \in \mathbb{R}^3 \mid |x| \leq k\}$ . Now, for  $u \in B$  and  $k > 0$ , we have:

$$\int_{\mathbb{R}^3} |(1 - \chi_k)\tilde{V}\tilde{u}|^p dx \leq \left( \int_{|x| \geq k} |\tilde{V}|^r dx \right)^{p/r} \left( \int_{|x| \geq k} |\tilde{u}|^{pr/(r-p)} dx \right)^{(r-p)/r}.$$

It follows that

$$|(1 - \chi_k)\tilde{V}\tilde{u}|_{L^p} \leq |B|_{L^{pr/(r-p)}} |(1 - \chi_k)\tilde{V}|_{L^r}, \quad u \in B, \quad k > 0.$$

Similarly, we have:

$$(4.2) \quad |\chi_k \tilde{V}\tilde{u}|_{L^p} \leq |\tilde{V}|_{L^r} |\chi_k \tilde{u}|_{L^{pr/(r-p)}}, \quad u \in H_0^1(\Omega), \quad k > 0.$$

Now, given  $\varepsilon > 0$ , we choose  $k > 0$  so large that  $|(1 - \chi_k)\tilde{V}|_{L^r} \leq \varepsilon$ . Then

$$\begin{aligned} \{\tilde{V}\tilde{u} \mid u \in B\} &= \{\chi_k \tilde{V}\tilde{u} + (1 - \chi_k)\tilde{V}\tilde{u} \mid u \in B\} \\ &\subset \{\chi_k \tilde{V}\tilde{u} \mid u \in B\} + \{(1 - \chi_k)\tilde{V}\tilde{u} \mid u \in B\} \\ &\subset \{v \in L^p(\mathbb{R}^3) \mid |v|_{L^p} \leq \varepsilon\} + \{\chi_k \tilde{V}\tilde{u} \mid u \in B\}. \end{aligned}$$

We notice that  $2 \leq pr/(r-p) < 6$ : therefore, By Rellich's Theorem,  $H^1(B_k(0))$  is compactly embedded into  $L^{pr/(r-p)}$ . It follows that the set  $\{\chi_k \tilde{u} \mid u \in B\}$  is precompact in  $L^{pr/(r-p)}$ . By (4.2), we deduce that  $\{\chi_k \tilde{V}\tilde{u} \mid u \in B\}$  is precompact in  $L^p(\mathbb{R}^3)$ . A simple *measure of non compactness argument* shows then that the set  $\{\tilde{V}\tilde{u} \mid u \in B\}$  is precompact in  $L^p(\mathbb{R}^3)$  and this in turn implies that the set  $\{Vu \mid u \in B\}$  is precompact in  $L^p(\Omega)$ .  $\square$

LEMMA 4.6. *Let  $V$  be as in Lemma 4.5. Let  $A+V$  be the selfadjoint operator determined by the bilinear form  $a(u, v) + \int_{\Omega} Vuv dx$ ,  $u, v \in H_0^1(\Omega)$ . Then, for sufficiently large  $\lambda > 0$ ,  $(A + \lambda)^{-1} - (A + V + \lambda)^{-1}$  is a compact operator in  $L^2(\Omega)$ .*

PROOF. Take  $\lambda > 0$  so large that  $A + V + \lambda$  be strictly positive. Let  $u \in L^2(\Omega)$ . Set  $v := (A + V + \lambda)^{-1}u$ ,  $w := (A + \lambda)^{-1}u$  and  $z := v - w$ . This means that

$$a(v, \phi) + \lambda(v, \phi) + \int_{\Omega} Vv\phi dx = \int_{\Omega} u\phi dx, \quad \text{for all } \phi \in H_0^1(\Omega)$$

and

$$a(w, \phi) + \lambda(w, \phi) = \int_{\Omega} u\phi dx, \quad \text{for all } \phi \in H_0^1(\Omega).$$

It follows that

$$a(z, \phi) + \lambda(z, \phi) + \int_{\Omega} Vv\phi dx = 0, \quad \text{for all } \phi \in H_0^1(\Omega).$$

Choosing  $\phi := z$ , Proposition 2.4 and Lemma 4.5 imply

$$\lambda_0 |z|_{H^1}^2 \leq |z|_{H^1} |Vv|_{H^{-1}} \leq \frac{\lambda_0}{2} |z|_{H^1}^2 + K_{\lambda_0} |Vv|_{H^{-1}}^2.$$

Therefore we obtain the estimate

$$|(A + \lambda)^{-1}u - (A + V + \lambda)^{-1}u|_{H^1} \leq K_{\lambda_0}|V(A + V + \lambda)^{-1}u|_{H^{-1}}, \quad u \in L^2(\Omega),$$

and the conclusion follows from Lemma 4.5. □

Now we can prove:

**PROPOSITION 4.7.** *Assume Hypothesis 4.4 is satisfied. Then the essential spectrum of  $A_\delta$  is contained in  $[(1 - \delta)\lambda_1, +\infty[$ .*

**PROOF.** Hypothesis 4.4 and Proposition 4.1 imply that, for every  $\varepsilon > 0$ , the bottom of the essential spectrum of  $A_\delta$  is larger than or equal to the bottom of the essential spectrum of  $(1 - \delta)A - \varepsilon - V_\varepsilon(x)$ . We observe that the spectrum of  $(1 - \delta)A - \varepsilon$  is contained in  $[(1 - \delta)\lambda_1 - \varepsilon, +\infty[$ . By Lemma 4.6 and Weyl's Theorem (see [19, Theorem XIII.14]), the essential spectrum of  $(1 - \delta)A - \varepsilon - V_\varepsilon(x)$  coincides with that of  $(1 - \delta)A - \varepsilon$ . It follows that the bottom of the essential spectrum of  $A_\delta$  is larger than or equal to  $(1 - \delta)\lambda_1 - \varepsilon$  for arbitrary small  $\varepsilon > 0$ , and the conclusion follows. □

Whenever Hypothesis 4.4 is satisfied, for  $0 < \delta < 1$  and  $\lambda < (1 - \delta)\lambda_0$  we introduce the following quantity:

$$\mathcal{N}(\delta, \lambda) := \# \text{ eigenvalues of } A_\delta \text{ below } \lambda.$$

Then, for  $d \geq \mathcal{N}(\delta, \frac{1-\delta}{2}\lambda_1)$  we have:

$$\sum_{i=1}^d \mu_i(A_\delta) \geq \mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right)\mu_1(\delta) + \left(d - \mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right)\right)\frac{1-\delta}{2}\lambda_1.$$

We have thus proved our first main result:

**THEOREM 4.8.** *Assume Hypotheses 2.2, 2.6 and 4.4 are satisfied. Let  $\mathcal{I} \subset H_0^1(\Omega)$  be a compact invariant set for the semiflow  $\pi$  generated by equation (2.2) in  $H_0^1(\Omega)$ . Then the Hausdorff dimension of  $\mathcal{I}$  in  $L^2(\Omega)$  is finite and less than or equal to  $d$ , provided  $d$  is an integer number larger than  $\max\{d_1, d_2\}$ , where*

$$d_1 := \mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right)$$

and

$$d_2 := \frac{2}{(1-\delta)\lambda_1} \left( \mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right) \left( \frac{1-\delta}{2}\lambda_1 - \mu_1(A_\delta) \right) + D(\gamma, \lambda_0, \delta, |\mathcal{I}|_{H^1}) \right).$$

**REMARK 4.9.** The first proper value  $\mu_1(A_\delta)$  of  $A_\delta$  can be estimated from below in terms of  $\lambda_0$  and  $|\partial_u f(\cdot, 0)|_{L_u^\sigma}$ . The explicit computations are left to the reader.

**REMARK 4.10.** By Lemma 2.8, also the Hausdorff dimension of  $\mathcal{I}$  in  $H_0^1(\Omega)$  is finite and it is equal to the Hausdorff dimension of  $\mathcal{I}$  in  $L^2(\Omega)$ .

### 5. Estimate of $\mathcal{N}(\delta, \frac{1-\delta}{2}\lambda_1)$

In this section we shall obtain an explicit estimate for the number  $\mathcal{N}(\delta, \frac{1-\delta}{2}\lambda_1)$  in terms of the dominating potential  $V_\varepsilon$  of Hypothesis 4.4. Our main tool is the celebrated Cwickel–Lieb–Rozenblum inequality, in its abstract formulation due to Rozenblum and Solomyak (see [21]). In order to exploit the CLR inequality, we need to make some assumption on the regularity of the open domain  $\Omega$ . Namely, we make the following assumption:

**HYPOTHESIS 5.1.** *The open set  $\Omega$  is a uniformly  $C^2$  domain in the sense of Browder [7, p. 36].*

As a consequence, by elliptic regularity we have that

$$D(-\Delta) = H^2(\Omega) \cap H_0^1(\Omega) \subset L^\infty(\Omega).$$

In this situation, if  $\omega \in L_u^\sigma(\mathbb{R}^3)$  then the assignment  $u \mapsto \omega u$  defines a relatively bounded perturbation of  $-\Delta$  and therefore  $D(-\Delta + \omega) = H^2(\Omega) \cap H_0^1(\Omega)$ . It follows that  $X^\alpha \subset L^\infty(\Omega)$  for  $\alpha > 3/4$  (see [11, Theorem 1.6.1]).

Set  $\bar{\varepsilon} := (1 - \delta)\lambda_1/4$ . Define the bilinear forms

$$\tilde{a}_{\delta, \bar{\varepsilon}}(u, v) := (1 - \delta) \left( \int_{\Omega} \nabla u \cdot \nabla v \, dx + \int_{\Omega} \beta uv \, dx \right) - 3\bar{\varepsilon} \int_{\Omega} uv \, dx,$$

for  $u, v \in H_0^1(\Omega)$ , and

$$b_{\delta, \bar{\varepsilon}}(u, v) := - \int_{\Omega} V_{\bar{\varepsilon}} uv \, dx.$$

Moreover, set

$$a_{\delta, \bar{\varepsilon}}(u, v) := \tilde{a}_{\delta, \bar{\varepsilon}}(u, v) + b_{\delta, \bar{\varepsilon}}(u, v)$$

and denote by  $\tilde{A}_{\delta, \bar{\varepsilon}}$  and  $A_{\delta, \bar{\varepsilon}}$  the selfadjoint operators induced by  $\tilde{a}_{\delta, \bar{\varepsilon}}$  and  $a_{\delta, \bar{\varepsilon}}$ , respectively.

A simple computation shows that

$$\mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right) \leq n_{\delta, \bar{\varepsilon}},$$

where  $n_{\delta, \bar{\varepsilon}}$  is the number of negative eigenvalues of  $A_{\delta, \bar{\varepsilon}}$ .

By Theorem 1.3.2 in [8], the operator  $\tilde{A}_{\delta, \bar{\varepsilon}}$  is positive (with  $\tilde{A}_{\delta, \bar{\varepsilon}} \geq \bar{\varepsilon}I$ ) and order preserving. Moreover, since  $D(A_{\delta, \bar{\varepsilon}}^\alpha) \subset L^\infty(\Omega)$  for  $\alpha > 3/4$ , then for every such  $\alpha$  and  $\gamma < \bar{\varepsilon}$  we have

$$|e^{-t\tilde{A}_{\delta, \bar{\varepsilon}}}u|_{L^\infty} \leq M_{\alpha, \gamma} t^{-\alpha} e^{-\gamma t} |u|_{L^2}, \quad u \in L^2(\Omega),$$

where  $M_{\alpha, \gamma}$  is a constant depending only on  $\alpha$ ,  $\gamma$  and on the embedding constant of  $H^2(\Omega)$  into  $L^\infty(\Omega)$ . It follows that

$$M_{\tilde{A}_{\delta, \bar{\varepsilon}}}^2(t) := \|e^{-(t/2)\tilde{A}_{\delta, \bar{\varepsilon}}}\|_{\mathcal{L}(L^2, L^\infty)}^2 \leq M_{\alpha, \gamma}^2 2^{2\alpha} t^{-2\alpha} e^{-\gamma t}.$$

We are now in a position to apply Theorem 2.1 in [21]. We have thus proved the following theorem:

**THEOREM 5.2.** *Assume that Hypotheses 2.2, 2.6, 4.4 and 5.1 are satisfied. Let  $\bar{\varepsilon} := (1 - \delta)\lambda_1/4$ . Then*

$$\mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right) \leq n_{\delta, \bar{\varepsilon}} \leq C_{q/2} M_{q/2, \gamma} \int_{\Omega} V_{\bar{\varepsilon}}(x)^q dx,$$

where  $C_{\alpha}$  is a constant depending only on  $\alpha$ , for  $\alpha > 3/4$ .

**6. Dissipative equations: dimension of the attractor**

In this section we specialize our results to the case of a dissipative equation. We make the following assumption:

**HYPOTHESIS 6.1.** *There exists a non negative function  $D \in L^q(\Omega)$ ,  $2 \leq q > 3/2$ , such that*

$$(6.1) \quad f(x, u)u \leq D(x)|u|, \quad (x, u) \in \Omega \times \mathbb{R}.$$

**REMARK 6.2.** Hypotheses 6.1 and 2.2 together are equivalent to the structure assumption of Theorem 4.4 in [4].

An easy computation shows that  $|f(x, 0)| \leq D(x)$  for  $x \in \Omega$ , and that  $F(x, u) := \int_0^u f(x, s) ds$  satisfies

$$F(x, u) \leq D(x)|u|, \quad (x, u) \in \Omega \times \mathbb{R}.$$

By slightly modifying some technical arguments in [17], one can prove that the semiflow  $\pi$  generated by equation (2.2) in  $H_0^1(\Omega)$  possesses a compact global attractor  $\mathcal{A}$ . Moreover,  $\pi$  is gradient-like with respect to the Lyapunov functional

$$\mathcal{L}(u) := \int_{\Omega} |\nabla u|^2 dx + \int_{\Omega} \beta(x)|u|^2 dx - \int_{\Omega} F(x, u) dx, \quad u \in H^1(\Omega).$$

Assuming Hypothesis 6.1, we shall give an explicit estimate for  $|\mathcal{A}|_{H^1}$  in terms of  $|D|_{L^q}$ . Moreover, we shall prove that Hypothesis 6.1 implies Hypothesis 4.4, and we explicitly compute the dominating potential  $V_{\bar{\varepsilon}}$  in terms of  $D$ . Therefore, we are able to obtain an explicit estimate for the number  $\mathcal{N}\left(\delta, \frac{1-\delta}{2}\lambda_1\right)$  in terms of  $|D|_{L^q}$ . As a consequence, the estimate of the dimension of  $\mathcal{A}$  given by Theorem 4.8 can be made completely explicit in terms of the structure parameters of equation (1.1).

We have the following theorem:

THEOREM 6.3. *Assume Hypotheses 2.2, 2.6 and 6.1 are satisfied.*

(a) *Let  $\phi \in H_0^1(\Omega)$  be an equilibrium of  $\pi$ . Then*

$$|\phi|_{H^1} \leq \frac{M_{q'}}{\lambda_0} |D|_{L^q},$$

*where  $M_{q'}$  is the embedding constant of  $H_0^1(\mathbb{R}^3)$  into  $L^{q'}(\mathbb{R}^3)$ .*

(b) *There exists a constant  $S > 0$  such that*

$$|u|_{H^1} \leq S \quad \text{for all } u \in \mathcal{A}.$$

*The constant  $S$  can be explicitly computed and depends only on  $C, \gamma, \sigma, \lambda_0, \Lambda_0, |D|_{L^q}, |\partial_u f(\cdot, 0)|_{L_{\mathbb{U}}^\sigma}$  and on the constants of Sobolev embeddings.*

PROOF. Let  $\phi \in H_0^1(\Omega)$  be an equilibrium of  $\pi$ . Then, for  $\varepsilon > 0$ , we have

$$\begin{aligned} \lambda_0 |\phi|_{H^1}^2 &\leq \int_{\Omega} |\nabla \phi|^2 dx + \int_{\Omega} \beta(x) |\phi|^2 dx = \int_{\Omega} f(x, \phi) \phi dx \leq \int_{\Omega} D(x) |\phi| dx \\ &\leq |D|_{L^q} |\phi|_{L^{q'}} \leq \varepsilon |\phi|_{L^{q'}}^2 + \frac{1}{4\varepsilon} |D|_{L^q}^2 \leq \varepsilon M_{q'}^2 |\phi|_{H^1}^2 + \frac{1}{4\varepsilon} |D|_{L^q}^2; \end{aligned}$$

choosing  $\varepsilon := \lambda_0 / (2M_{q'}^2)$  we get property (a). In order to prove (b), we notice that, since  $\mathcal{L}$  is a Lyapunov functional for  $\pi$  and  $\mathcal{A}$  is compact in  $H_0^1(\Omega)$ , there exists an equilibrium  $\phi$  such that, for every  $u \in \mathcal{A}$ ,

$$\begin{aligned} \int_{\Omega} |\nabla u|^2 dx + \int_{\Omega} \beta(x) |u|^2 dx - \int_{\Omega} F(x, u) dx \\ \leq \int_{\Omega} |\nabla \phi|^2 dx + \int_{\Omega} \beta(x) |\phi|^2 dx - \int_{\Omega} F(x, \phi) dx. \end{aligned}$$

Then, for  $\varepsilon > 0$ , we have:

$$\begin{aligned} \lambda_0 |u|_{H^1}^2 &\leq \int_{\Omega} D(x) |u| dx + \Lambda_0 |\phi|_{H^1}^2 + \int_{\Omega} F(x, \phi) dx \\ &\leq \varepsilon M_{q'}^2 |u|_{H^1}^2 + \frac{1}{4\varepsilon} |D|_{L^q}^2 + \Lambda_0 |\phi|_{H^1}^2 + \int_{\Omega} F(x, \phi) dx. \end{aligned}$$

We choose  $\varepsilon := \lambda_0 / (2M_{q'}^2)$  and the conclusion follows.  $\square$

Finally, we have:

THEOREM 6.4. *Assume that Hypotheses 2.6 and 6.1 are satisfied. Then for every  $0 < \varepsilon \leq 1$ ,*

$$\partial_u f(x, 0) \leq \frac{2}{\varepsilon} D(x) + \frac{\varepsilon}{2} C(1 + \varepsilon^\gamma).$$

PROOF. For  $\varepsilon > 0$  we have:

$$f(x, \varepsilon) = f(x, 0) + \partial_u f(x, 0) \varepsilon + \int_0^\varepsilon \left( \int_0^s \partial_{uu} f(x, r) dr \right) ds.$$

It follows that

$$f(x, 0)\varepsilon + \partial_u f(x, 0)\varepsilon^2 + \varepsilon \int_0^\varepsilon \left( \int_0^s \partial_{uu} f(x, r) dr \right) ds = f(x, \varepsilon)\varepsilon \leq D(x)\varepsilon.$$

Therefore

$$\partial_u f(x, 0) \leq \frac{D(x) + |f(x, 0)|}{\varepsilon} + \frac{1}{\varepsilon} \int_0^\varepsilon \left( \int_0^s C(1 + |r|^\gamma) dr \right) ds,$$

and the conclusion follows.  $\square$

REMARK 6.5. Theorem 6.4 shows that Hypotheses 2.6 and 6.1 together imply Hypothesis 4.4, with  $V_\varepsilon(x) = \frac{2C}{\varepsilon} D(x)$ .

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MARTINO PRIZZI  
Università di Trieste  
Dipartimento di Matematica e Informatica  
Via Valerio 12/1  
34127 Trieste, ITALY  
*E-mail address:* mprizzi@units.it