# Pullback attractors for semilinear non-autonomous problem 

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## The problem

Let $\Omega$ be a bounded domain in $\mathbb{R}^{N}$ with $N \geqslant 3$, which boundary $\partial \Omega$ is sufficiently regular. We consider the following initialboundary value problem

$$
\left\{\begin{array}{l}
u_{t t}-\Delta u+u+\eta(-\Delta)^{\frac{1}{2}} u_{t}+a_{\epsilon}(t)(-\Delta)^{\frac{1}{2}} v_{t}=f(u),  \tag{1}\\
v_{t t}-\Delta v+\eta(-\Delta)^{\frac{1}{2}} v_{t}-a_{\epsilon}(t)(-\Delta)^{\frac{1}{2}} u_{t}=0
\end{array}\right.
$$

$(x, t) \in \Omega \times(\tau, \infty)$, where $\eta$ is a positive constant, subject to boundary conditions

$$
\begin{equation*}
u=v=0,(x, t) \in \partial \Omega \times(\tau, \infty) \tag{2}
\end{equation*}
$$

and initial conditions

$$
\begin{array}{r}
u(\tau, x)=u_{0}(x), u_{t}(\tau, x)=u_{1}(x) \\
v(\tau, x)=v_{0}(x), v_{t}(\tau, x)=v_{1}(x), x \in \Omega, \tau \in \mathbb{R} \tag{3}
\end{array}
$$

Assume that the function $a_{\epsilon}: \mathbb{R} \rightarrow(0, \infty)$ is continuously differentiable in $\mathbb{R}$ and satisfies the following condition:

$$
\begin{equation*}
0<a_{0} \leq a_{\epsilon}(t) \leq a_{1}, \tag{4}
\end{equation*}
$$

for all $\epsilon \in[0,1]$ and $t \in \mathbb{R}$, with positive constants $a_{0}$ and $a_{1}$, and we also assume that the first derivative of $a_{\epsilon}$ is uniformly bounded in $t$ and $\epsilon$, that is, there exists a constant $b_{0}>0$ such that

$$
\begin{equation*}
\left|a_{\epsilon}^{\prime}(t)\right| \leq b_{0} \quad \text { for all } \quad t \in \mathbb{R}, \epsilon \in[0,1] . \tag{5}
\end{equation*}
$$

Furthermore, we assume that $a_{\epsilon}$ is $(\beta, C)$-Hölder continuous, for each $\epsilon \in[0,1]$; that is,

$$
\begin{equation*}
\left|a_{\epsilon}(t)-a_{\epsilon}(s)\right| \leq C|t-s|^{\beta} \tag{6}
\end{equation*}
$$

for all $t, s \in \mathbb{R}$ and $\epsilon \in[0,1]$.

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$$

for all $t, s \in \mathbb{R}$ and $\epsilon \in[0,1]$. Concerning the nonlinearity $f$, we assume that $f \in C^{1}(\mathbb{R})$ and it satisfies the dissipativeness condition

$$
\begin{equation*}
\limsup _{|s| \rightarrow \infty} \frac{f(s)}{s} \leq 0 \tag{7}
\end{equation*}
$$

and also satisfies the subcritical growth condition given by

$$
\begin{equation*}
\left|f^{\prime}(s)\right| \leq c\left(1+|s|^{\rho-1}\right), \tag{8}
\end{equation*}
$$

for all $s \in \mathbb{R}$, where $1<\rho<\frac{n}{n-2}$, with $n \geq 3$, and $c>0$ is a constant.

In the case that $a_{\epsilon}(t) \equiv a$, the system (1) represents the autonomous version of the Klein-Gordon-Zakharov system. Within the autonomous case, if $n=3$ then the Klein-Gordon-Zakharov system arises to describe the interaction of a Langmuir wave (Plasma oscillations, are rapid oscillations of the electron density in conducting media such as plasmas or metals in the ultraviolet region) and acoustic wave in a plasma.

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These types of systems have been considered by many researchers in recent years.

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For more details see [Bonotto, Nascimento and Santiago, Long-time behaviour for a non-autonomous Klein-Gordon-Zakharov system, Journal of Mathematical Analysis and Applications, 506 (2022), 125670.

## Basic Concepts

Suppose that we have a non-autonomous differential equations in a Banach space $X$

$$
\frac{d u}{d t}=f(t, u), \quad u(s)=u_{0}
$$

with a unique solution $u\left(t, s, u_{0}\right)$. Note that the initial time has a very important role because we have an explicit dependence on time of $f$. This time dependence may appear in external force, in the operator, in both at the same time or even on the boundary conditions.

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In general, a non-autonomous system shows two different important dynamics without relation between them:
(1) forward dynamic: the behavior when final time goes to infinity:

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(1) forward dynamic: the behavior when final time goes to infinity:

$$
\lim _{t \rightarrow \infty} u\left(t, s, u_{0}\right) .
$$

(2) pullback dynamics: the behavior when the initial time goes to minus infinity:

$$
\lim _{s \rightarrow-\infty} u\left(t, s, u_{0}\right)
$$

| Semigroup | Process |
| :---: | :---: |
| $T(t)$ | $U(t, s)$ |
| exponential decay | exponential stability |
| $\\|T(t)\\| \leqslant e^{-\beta t}$ | $\\|U(t, s)\\| \leqslant e^{-\beta(t-s)}$ |
| Invariance | Invariance |
| $T(t) A=A$ | $U(t, s) A(s)=A(t)$ |
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For more details see [Carvalho, Langa and Robinson, Attractors for infinite-dimensional non-autonomous semantical systems, Springer, 2012.]

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## Singularly non-autonomous abstract problem

Here, $L(\mathcal{Z})$ will denote the space of linear and bounded operators defined in a Banach space $\mathcal{Z}$. Let $\mathcal{A}(t), t \in \mathbb{R}$, be a family of unbounded closed linear operators defined on a fixed dense subspace $D$ of $\mathcal{Z}$.
Consider the singularly non-autonomous parabolic problem

$$
\left\{\begin{array}{l}
\frac{d u}{d t}+\mathcal{A}(t) u=0, t>\tau  \tag{9}\\
u(\tau)=I
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$$

We assume
(a) The operator $\mathcal{A}(t): D \subset \mathcal{Z} \rightarrow \mathcal{Z}$ is a closed densely defined operator (the domain $D$ is fixed) and there is a constant $C>0$ (independent of $t \in \mathbb{R}$ ) such that
$\left\|(\lambda I+\mathcal{A}(t))^{-1}\right\|_{L(\mathcal{Z})} \leqslant \frac{C}{|\lambda|+1} ;$ for all $\lambda \in \mathbb{C}$ with $\operatorname{Re} \lambda \geqslant 0$.
To express this fact we will say that the family $\mathcal{A}(t)$ is uniformly sectorial.
(b) There are constants $C>0$ and $\epsilon_{0}>0$ such that, for any $t, \tau, s \in \mathbb{R}$,

$$
\left\|[\mathcal{A}(t)-\mathcal{A}(\tau)] \mathcal{A}^{-1}(s)\right\|_{L(\mathcal{Z})} \leqslant C(t-\tau)^{\epsilon_{0}}, \quad \epsilon_{0} \in(0,1] .
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To express this fact we will say that the map $\mathbb{R} \ni t \mapsto \mathcal{A}(t)$ is uniformly Hölder continuous.
Denote by $\mathcal{A}_{0}$ the operator $\mathcal{A}\left(t_{0}\right)$ for some $t_{0} \in \mathbb{R}$ fixed. If $\mathcal{Z}^{\alpha}$ denotes the domain of $\mathcal{A}_{0}^{\alpha}, \alpha>0$, with the graph norm and $\mathcal{Z}^{0}:=$ $\mathcal{Z}$, denote by $\left\{\mathcal{Z}^{\alpha} ; \alpha \geqslant 0\right\}$ the fractional power scale associated with $\mathcal{A}_{0}$ (see Henry [Springer, 1981] and Amann [Birkhäuser Verlag, Basel, 1995].

From $(a),-\mathcal{A}(t)$ is the infinitesimal generator of an analytic semigroup $\left\{e^{-\tau \mathcal{A}(t)} \in L(\mathcal{Z}): \tau \geqslant 0\right\}$. Using this and the fact that $0 \in \rho(\mathcal{A}(t))$, it follows that

$$
\left\|e^{-\tau \mathcal{A}(t)}\right\|_{L(\mathcal{Z})} \leqslant C e^{-\delta \tau}, \delta>0, \tau \geqslant 0, t \in \mathbb{R}
$$

and

$$
\left\|\mathcal{A}(t) e^{-\tau \mathcal{A}(t)}\right\|_{L(\mathcal{Z})} \leqslant C \tau^{-1} e^{-\delta \tau}, \delta>0, \tau>0, t \in \mathbb{R}
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$$

It follows from $(b)$ that $\left\|\mathcal{A}(t) \mathcal{A}^{-1}(\tau)\right\|_{L(\mathcal{Z})} \leqslant C, \forall t, \tau \in I$, for all $I \subset \mathbb{R}$ bounded. Also, the semigroup $e^{-\tau \mathcal{A}(t)}$ generated by $-\mathcal{A}(t)$ satisfies the following estimate

$$
\left\|e^{-\tau \mathcal{A}(t)}\right\|_{L\left(\mathcal{Z}^{\beta}, \mathcal{Z}^{\alpha}\right)} \leqslant M \tau^{\beta-\alpha}
$$

where $0 \leqslant \beta \leqslant \alpha<1+\epsilon_{0}$.

Next we recall the definition of a linear evolution process associated with a family of operators $\{\mathcal{A}(t): t \in \mathbb{R}\}$.

## Definition 1

A family $\{L(t, \tau): t \geqslant \tau \in \mathbb{R}\} \subset L(\mathcal{Z})$ satisfying

$$
\begin{aligned}
& \text { 1) } L(\tau, \tau)=I \\
& \text { 2) } L(t, \sigma) L(\sigma, \tau)=L(t, \tau) \text {, for any } t \geqslant \sigma \geqslant \tau \\
& \text { 3) } \mathcal{P} \times \mathcal{Z} \ni\left((t, \tau), u_{0}\right) \mapsto L(t, \tau) v_{0} \in \mathcal{Z} \text { is continuous, }
\end{aligned}
$$

where $\mathcal{P}=\left\{(t, \tau) \in \mathbb{R}^{2}: t \geqslant \tau\right\}$, is called a linear evolution process (process for short) or family of evolution operators.

If the operator $\mathcal{A}(t)$ is uniformly sectorial and uniformly Hölder continuous, then there exists a linear evolution process $\{L(t, \tau)$ : $t \geqslant \tau \in \mathbb{R}\}$ associated with $\mathcal{A}(t)$, which is given by

$$
L(t, \tau)=e^{-(t-\tau) \mathcal{A}(\tau)}+\int_{\tau}^{t} L(t, s)[\mathcal{A}(\tau)-\mathcal{A}(s)] e^{-(s-\tau) \mathcal{A}(\tau)} d s
$$

that is solution of (9).
For more details see Carvalho and Nascimento [DCDS, 2009].

We consider the singularly non-autonomous abstract parabolic problem

$$
\left\{\begin{array}{l}
\frac{d u}{d t}+\mathcal{A}(t) u=g(t, u), t>\tau  \tag{10}\\
u(\tau)=u_{0} \in D
\end{array}\right.
$$

where the operator $\mathcal{A}(t)$ is uniformly sectorial and uniformly Hölder continuous and the nonlinearity $g$ satisfies conditions which will be specified later.

## Definition 2

Let $g: \mathbb{R} \times X^{\alpha} \rightarrow X^{\beta}, \alpha \in[\beta, \beta+1)$ be a continuous function. We say that a function $u$ is a (local) solution of (10) starting in $u_{0} \in X^{\alpha}$, if $u \in C\left(\left[\tau, \tau+t_{0}\right], X^{\alpha}\right) \cap C^{1}\left(\left(\tau, \tau+t_{0}\right], X^{\alpha}\right), u(\tau)=u_{0}$, $u(t) \in D(\mathcal{A}(t))$ for all $t \in\left(\tau, \tau+t_{0}\right]$ and (10) is satisfied for all $t \in\left(\tau, \tau+t_{0}\right)$.

Now we state the following abstract local well-posedness result.

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Now we state the following abstract local well-posedness result.

## Theorem 3 (Caraballo et al. Nonlinear Analysis (2010))

Suppose that the family of operators $\mathcal{A}(t)$ is uniformly sectorial and uniformly Hölder continuous in $X^{\beta}$. If $g: \mathbb{R} \times X^{\alpha} \rightarrow X^{\beta}$, $\alpha \in[\beta, \beta+1)$, is a Lipschitz continuous map in bounded subsets of $X^{\alpha}$, then, given $r>0$, there is a time $t_{0}>0$ such that for all $u_{0} \in B_{X^{\alpha}}(0 ; r)$ there exists a unique solution $u\left(\cdot, \tau, u_{0}\right) \in$ $C\left(\left[\tau, \tau+t_{0}\right], X^{\alpha}\right) \cap C^{1}\left(\left(\tau, \tau+t_{0}\right], X^{\alpha}\right)$ of the problem (10) starting in $u_{0} \in X^{\alpha}$. Moreover, such solutions are continuous with respect the initial data in $B_{X^{\alpha}}(0 ; r)$.

## Basic definitions and existence results

We start remembering the definition of Hausdorff semi-distance between two subsets $A$ and $B$ of a metric space $(X, d)$ :

$$
\operatorname{dist}_{H}(A, B)=\sup _{a \in A} \inf _{b \in B} d(a, b) .
$$

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$$

Next we present several definitions about theory of pullback attractors.

## Definition 4

Let $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ be an evolution process in a metric space $X$. A set $B(t) \subset X$ pullback attracts a set $C$ at time $t$ under $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ if

$$
\lim _{\tau \rightarrow-\infty} \operatorname{dist}_{X}(S(t, \tau) C, B(t))=0
$$

where $S(t, \tau) C:=\{S(t, \tau) x \in X: x \in C\}$.

## Definition 5

We say that an evolution process $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ in $X$ is pullback strongly bounded if, for each $t \in \mathbb{R}$ and each bounded subset $B$ of $X$,

$$
\bigcup_{\tau \leqslant t} \bigcup_{s \leqslant \tau} S(\tau, s) B
$$

is bounded.

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## Definition 6

An evolution process $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ in $X$ is pullback asymptotically compact if, for each $t \in \mathbb{R}$, each sequence $\left\{\tau_{n}\right\}$ in $(-\infty, t]$ with $\tau_{n} \rightarrow-\infty$ as $n \rightarrow \infty$ and each bounded sequence $\left\{x_{n}\right\}$ in $X$ such that $\left\{S\left(t, \tau_{n}\right) x_{n}\right\} \subset X$ is bounded, the sequence $\left\{S\left(t, \tau_{n}\right) x_{n}\right\}$ is relatively compact in $X$.

## Definition 7

A family $\{\mathbb{A}(t): t \in \mathbb{R}\}$ of compact subsets of $X$ is a pullback attractor for an evolution process $\{S(t, \tau): t \geq \tau \in \mathbb{R}\}$ if the following conditions hold:
(i) $\{\mathbb{A}(t): t \in \mathbb{R}\}$ is invariant, that is, $S(t, \tau) \mathbb{A}(\tau)=\mathbb{A}(t)$ for all $t \geq \tau$,
(ii) $\{\mathbb{A}(t): t \in \mathbb{R}\}$ pullback attracts bounded subsets of $X$, that is,

$$
\lim _{\tau \rightarrow-\infty} \mathrm{d}_{\mathrm{X}}(S(t, \tau) B, \mathbb{A}(t))=0
$$

for every $t \in \mathbb{R}$ and every bounded subset $B$ of $X$, where $S(t, \tau) B=\{S(t, \tau) x: x \in B\}$ is the image of $B$ under $\{S(t, \tau): t \geq \tau \in \mathbb{R}\}$, and
(iii) $\{\mathbb{A}(t): t \in \mathbb{R}\}$ is the minimal family of closed sets satisfying property (ii).

In applications, to prove that a process has a pullback attractor we use the Theorem below, proved in Caraballo et al., [Nonlinear Anal. (2010)] which gives a sufficient condition for existence of a pullback attractor.

## Definition 8

An evolution process $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ in $X$ is pullback strongly bounded dissipative if, for each $t \in \mathbb{R}$, there is a bounded subset $B(t)$ of $X$ which pullback absorbs bounded subsets of $X$ at time $s$ for each $s \leqslant t$; that is, given a bounded subset $B$ of $X$ and $s \leqslant t$, there exists $\tau_{0}(s, B)$ such that $S(s, \tau) B \subset B(t)$, for all $\tau \leqslant \tau_{0}(s, B)$.

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## Theorem 9 (CCLR, 2010)

If an evolution process $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ in the metric space $X$ is pullback strongly bounded dissipative and pullback asymptotically compact, then $\{S(t, \tau): t \geqslant \tau \in \mathbb{R}\}$ has a pullback attractor $\{\mathbb{A}(t): t \in \mathbb{R}\}$ with the property that $\cup_{\tau \leqslant t} \mathbb{A}(\tau)$ is bounded for each $t \in \mathbb{R}$.

## Theorem 10 (CCLR, 2010)

Let $\{S(t, s): t \geqslant s \in \mathbb{R}\}$ be a pullback strongly bounded evolution process such that $S(t, s)=L(t, s)+U(t, s)$, where there exist a non-increasing function $k: \mathbb{R}^{+} \times \mathbb{R}^{+} \rightarrow \mathbb{R}$, with $k(\sigma, r) \rightarrow 0$ when $\sigma \rightarrow \infty$, and for all $s \leqslant t$ and $x \in X$ with $\|x\| \leqslant r,\|L(t, s) x\| \leqslant$ $k(t-s, r)$, and $U(t, s)$ is compact. Then, the family of evolution process $\{S(t, s): t \geqslant s \in \mathbb{R}\}$ is pullback asymptotically compact.

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For more details see [Carvalho, Langa and Robinson, Attractors for infinite-dimensional non-autonomous semantical systems, Springer, 2012.]

## Abstract setting

Consider the following initial-boundary value problem

$$
\left\{\begin{array}{l}
u_{t t}-\Delta u+u+\eta(-\Delta)^{\frac{1}{2}} u_{t}+a_{\epsilon}(t)(-\Delta)^{\frac{1}{2}} v_{t}=f(u),  \tag{11}\\
v_{t t}-\Delta v+\eta(-\Delta)^{\frac{1}{2}} v_{t}-a_{\epsilon}(t)(-\Delta)^{\frac{1}{2}} u_{t}=0
\end{array}\right.
$$

$(x, t) \in \Omega \times(\tau, \infty)$, where $\eta$ is a positive constant, subject to boundary conditions

$$
\begin{equation*}
u=v=0,(x, t) \in \partial \Omega \times(\tau, \infty) \tag{12}
\end{equation*}
$$

and initial conditions

$$
\begin{array}{r}
u(\tau, x)=u_{0}(x), u_{t}(\tau, x)=u_{1}(x),  \tag{13}\\
v(\tau, x)=v_{0}(x), v_{t}(\tau, x)=v_{1}(x), x \in \Omega, \tau \in \mathbb{R}
\end{array}
$$

In order to formulate the non-autonomous problem (1) - (3) in a nonlinear evolution process setting, we introduce some notations. Let $X=L^{2}(\Omega)$ and denote by $A: D(A) \subset X \rightarrow X$ the negative Laplacian operator, that is, $A u=(-\Delta) u$ for all $u \in D(A)$, where $D(A)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Thus $A$ is a positive self-adjoint operator in $X$ with compact resolvent and, therefore, $-A$ generates a compact analytic semigroup on $X$. Following Henry [Springer, 1981]), $A$ is a sectorial operator in $X$. Now, denote by $X^{\alpha}$, $\alpha>0$, the fractional power spaces associated with the operator $A$; that is, $X^{\alpha}=D\left(A^{\alpha}\right)$ endowed with the graph norm. With this notation, we have $X^{-\alpha}=\left(X^{\alpha}\right)^{\prime}$ for all $\alpha>0$, see Amann (Birkhäuser Verlag, 1995).

In this framework, the non-autonomous problem (1) - (3) can be rewritten as an ordinary differential equation in the following abstract form

$$
\begin{cases}W_{t}+\mathcal{A}(t) W=F(W), & t>\tau  \tag{14}\\ W(\tau)=W_{0}, & \tau \in \mathbb{R}\end{cases}
$$

where $W=W(t)$, for all $t \in \mathbb{R}$, and $W_{0}=W(\tau)$ are respectively given by

$$
W=\left[\begin{array}{c}
u \\
u_{t} \\
v \\
v_{t}
\end{array}\right] \text { and } \quad W_{0}=\left[\begin{array}{c}
u_{0} \\
u_{1} \\
v_{0} \\
v_{1}
\end{array}\right]
$$

and, for each $t \in \mathbb{R}$, the unbounded linear operator $\mathcal{A}(t): D(\mathcal{A}(t)) \subset$ $Y \rightarrow Y$ is defined by

$$
\mathcal{A}(t)\left[\begin{array}{c}
u  \tag{15}\\
v \\
w \\
z
\end{array}\right]=\left[\begin{array}{c}
-v \\
(A+I) u+\eta A^{\frac{1}{2}} v+a_{\epsilon}(t) A^{\frac{1}{2}} z \\
-z \\
-a_{\epsilon}(t) A^{\frac{1}{2}} v+A w+\eta A^{\frac{1}{2}} z
\end{array}\right]
$$

for each $\left[\begin{array}{llll}u & v & w & z\end{array}\right]^{T}$ in the domain $D(\mathcal{A}(t))$ defined by the space

$$
\begin{equation*}
D(\mathcal{A}(t))=Y^{1}=X^{1} \times X^{\frac{1}{2}} \times X^{1} \times X^{\frac{1}{2}}, \tag{16}
\end{equation*}
$$

where

$$
Y=Y_{0}=X^{\frac{1}{2}} \times X \times X^{\frac{1}{2}} \times X
$$

is the phase space of the problem (1) - (3).

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$$
Y=Y_{0}=X^{\frac{1}{2}} \times X \times X^{\frac{1}{2}} \times X
$$

is the phase space of the problem (1) $-(3)$. The nonlinearity $F$ is given by

$$
F(W)=\left[\begin{array}{c}
0  \tag{17}\\
f^{e}(u) \\
0 \\
0
\end{array}\right],
$$

where $f^{e}(u)$ is the Nemitskii operator associated with $f(u)$; that is,

$$
f^{e}(u)(x)=f(u(x)), \quad \text { for all } \quad x \in \Omega .
$$

Now, we observe that the norms

$$
\|(x, y, z, w)\|_{1}=\|x\|_{X^{\frac{1}{2}}}+\|y\|_{X}+\|z\|_{X^{\frac{1}{2}}}+\|w\|_{X}
$$

and

$$
\|(x, y, z, w)\|_{2}=\left(\|x\|_{X^{\frac{1}{2}}}^{2}+\|y\|_{X}^{2}+\|z\|_{X^{\frac{1}{2}}}^{2}+\|w\|_{X}^{2}\right)^{\frac{1}{2}}
$$

are equivalent in $Y_{0}$. In this way, we shall use the same notation $\|(x, y, z, w)\|_{Y_{0}}$ for both norms and the choice will be as convenient.

In the abstract form our problem was written as

$$
W_{t}+\mathcal{A}(t) W=F(W)
$$

where

$$
\mathcal{A}(t): Y^{1} \rightarrow Y
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$$

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$$

However, the natural energy of the problem is given by

$$
\begin{aligned}
\mathcal{E}(t) & =\frac{1}{2}\|u(t)\|_{X^{\frac{1}{2}}}^{2}+\frac{1}{2}\|u(t)\|_{X}^{2}+\frac{1}{2}\left\|u_{t}(t)\right\|_{X}^{2}+\frac{1}{2}\|v(t)\|_{X^{\frac{1}{2}}}^{2} \\
& +\frac{1}{2}\left\|v_{t}(t)\right\|_{X}^{2}-\int_{\Omega} \int_{0}^{u} f(s) d s d x
\end{aligned}
$$

that is naturally defined in $Y$ and we will show that $\mathcal{E}(t)$ decays along the solutions,

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\end{aligned}
$$

that is naturally defined in $Y$ and we will show that $\mathcal{E}(t)$ decays along the solutions, and then we would like to write our problem in the form

$$
\mathcal{A}(t): Y \rightarrow Y_{-1}
$$

for some appropriate space $Y_{-1}$ such that $Y \subset Y_{-1}$ and

$$
F(\cdot): Y \rightarrow Y_{\alpha-1} \hookrightarrow Y_{-1}, \quad(0<\alpha<1)
$$





## Theorem 11 (Well-Posedness)

Let $f \in C^{1}(\mathbb{R})$ be a function satisfying (7)-(8), assume conditions (4)-(6) hold and let $F: Y_{0} \rightarrow Y_{\alpha-1} \subset Y_{-1}$ be defined in (17). Then for any initial data $W_{0} \in Y_{0}$ the problem (14) has a unique global solution $W(t)$ such that

$$
W(t) \in C\left([\tau, \infty), Y_{0}\right)
$$

Moreover, such solutions are continuous with respect to the initial data on $Y_{0}$. Here, $Y_{-1}=X \times X^{-\frac{1}{2}} \times X \times X^{-\frac{1}{2}}$.

## Theorem 12 (Pullback Attractors)

Under the conditions of Theorem 11, the problem (1) - (3) has a pullback attractor $\{\mathbb{A}(t): t \in \mathbb{R}\}$ in $Y_{0}$ and

$$
\bigcup_{t \in \mathbb{R}} \mathbb{A}(t) \subset Y_{0}
$$

is bounded.

## Linear Analysis

Consider the linear problem associated with (1)-(3), in this case we consider the singularly non-autonomous linear parabolic problem

$$
\left\{\begin{array}{l}
w_{t}+\mathcal{A}(t) w=0, t>\tau  \tag{18}\\
w(\tau)=I
\end{array}\right.
$$

It is not difficult to see that $\operatorname{det}(\mathcal{A}(t))=A(A+I)$, and therefore that $0 \in$ $\rho(\mathcal{A}(t))$, for all $t \in \mathbb{R}$.

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$$

It is not difficult to see that $\operatorname{det}(\mathcal{A}(t))=A(A+I)$, and therefore that $0 \in$ $\rho(\mathcal{A}(t))$, for all $t \in \mathbb{R}$. Moreover, for each $t \in \mathbb{R}$, the operator $\mathcal{A}^{-1}(t): Y_{0} \rightarrow$ $Y_{0}$ is defined by
$\mathcal{A}^{-1}(t)\left[\begin{array}{c}u \\ v \\ w \\ z\end{array}\right]=\left[\begin{array}{cccc}\eta A^{\frac{1}{2}}(A+I)^{-1} & (A+I)^{-1} & a_{\epsilon}(t) A^{\frac{1}{2}}(A+I)^{-1} & 0 \\ -I & 0 & 0 & 0 \\ -a_{\epsilon}(t) A^{-\frac{1}{2}} & 0 & \eta A^{-\frac{1}{2}} & A^{-1} \\ 0 & 0 & -I & 0\end{array}\right]\left[\begin{array}{c}u \\ v \\ w \\ z\end{array}\right]$.

## Proposition 13

If $Y_{-1}$ denotes the extrapolation space of $Y_{0}=X^{\frac{1}{2}} \times X \times X^{\frac{1}{2}} \times X$ generated by the operator $\mathcal{A}^{-1}(t)$, then

$$
Y_{-1}=X \times X^{-\frac{1}{2}} \times X \times X^{-\frac{1}{2}}
$$

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$$
Y_{-1}=X \times X^{-\frac{1}{2}} \times X \times X^{-\frac{1}{2}}
$$

Recall that the extrapolation space $Y_{-1}$ is the completion of the normed space $\left(Y,\left\|\mathcal{A}^{-1}(t) \cdot\right\|_{Y}\right)$.

## Proposition 14

The family of operators $\{\mathcal{A}(t): t \in \mathbb{R}\}$, defined in (15) - (16), is uniformly Hölder continuous in $Y_{-1}$.

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## Proposition 14

The family of operators $\{\mathcal{A}(t): t \in \mathbb{R}\}$, defined in (15) - (16), is uniformly Hölder continuous in $Y_{-1}$.

Proof: Using (6), this result follows immediately from (15) and (16).

The next step is to show the analyticity of the semigroup $\left\{e^{-\tau \mathcal{A}(t)}\right.$ : $\tau \geq 0\}$.

## Theorem 15

The semigroup $\left\{e^{-\tau \mathcal{A}(t)}: \tau \geq 0\right\}$, generated by $-\mathcal{A}(t)$, is analytic for each $t \in \mathbb{R}$.

The next step is to show the analyticity of the semigroup $\left\{e^{-\tau \mathcal{A}(t)}\right.$ : $\tau \geq 0\}$.

## Theorem 15

The semigroup $\left\{e^{-\tau \mathcal{A}(t)}: \tau \geq 0\right\}$, generated by $-\mathcal{A}(t)$, is analytic for each $t \in \mathbb{R}$.

For the proof see [Bonotto, Nascimento and Santiago, JMAA (2022)].

## Theorem 16

The operators $\mathcal{A}(t)$ are uniformly sectorial and the map $\mathbb{R} \ni t \mapsto$ $\mathcal{A}(t) \in \mathcal{L}\left(Y_{-1}\right)$ is uniformly Hölder continuous. Then, for each functional parameter $a$, there exist a process

$$
\{L(t, \tau): t \geqslant \tau \in \mathbb{R}\}
$$

(or simply $L(t, \tau)$ ) associated with the operator $\mathcal{A}(t)$, that is solution of the linear problem associated with (1)-(3).

## Remark 17

We have the following description of the fractional power scale for the operator $\mathcal{A}(t)$, given as follows

$$
Y_{0} \hookrightarrow Y_{\alpha-1} \hookrightarrow Y_{-1}, \quad \text { for all } \quad 0<\alpha<1,
$$

where

$$
\begin{aligned}
Y_{\alpha-1} & =\left[Y_{-1}, Y_{0}\right]_{\alpha}=\left[X \times X^{-\frac{1}{2}} \times X \times X^{-\frac{1}{2}}, X^{\frac{1}{2}} \times X \times X^{\frac{1}{2}} \times X\right]_{\alpha} \\
& =\left[X, X^{\frac{1}{2}}\right]_{\alpha} \times\left[X^{-\frac{1}{2}}, X\right]_{\alpha} \times\left[X, X^{\frac{1}{2}}\right]_{\alpha} \times\left[X^{-\frac{1}{2}}, X\right]_{\alpha} \\
& =X^{\frac{\alpha}{2}} \times X^{\frac{\alpha-1}{2}} \times X^{\frac{\alpha}{2}} \times X^{\frac{\alpha-1}{2}},
\end{aligned}
$$

where $[\cdot, \cdot]_{\alpha}$ denotes the complex interpolation functor (see Triebel (1978)). The first equality follows from Proposition ?? (since $\left.0 \in \rho\left(A_{(a)}(t)\right)\right)$ see Amann (Example 4.7.3 (b)]) and the second equality follows from Proposition 2 in Carvalho and Cholewa (Bull. Austral. Math. Soc., (2002).

## Existence of local solutions

Proposition 18 gives us sufficient conditions for $F: Y_{0} \rightarrow Y_{\alpha-1}$ to be Lipschitz continuous in bounded subsets of $Y_{0}$.

## Proposition 18

Assume that $1<\rho<\frac{n+2(1-\alpha)}{n-2}$, with $\alpha \in(0,1)$. Then the map $F: Y_{0} \rightarrow Y_{\alpha-1}$, defined in (17), is Lipschitz continuous in bounded subsets of $Y_{0}$.

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## Proposition 18

Assume that $1<\rho<\frac{n+2(1-\alpha)}{n-2}$, with $\alpha \in(0,1)$. Then the map $F: Y_{0} \rightarrow Y_{\alpha-1}$, defined in (17), is Lipschitz continuous in bounded subsets of $Y_{0}$.

## Corollary 19

Let $1<\rho<\frac{n+2(1-\alpha)}{n-2}$, with $\alpha \in(0,1), f \in C^{1}(\mathbb{R})$ be a function satisfying (7)-(8), assume conditions (4)-(6) hold and let $F: Y_{0} \rightarrow Y_{\alpha-1}$ be defined in (17). Then given $r>0$, there exists a time $t_{0}=t_{0}(r)>0$ such that for all $W_{0} \in B_{Y_{0}}(0, r)$, there exists a unique solution $W:\left[\tau, \tau+t_{0}\right] \rightarrow Y_{0}$ of the problem (14) starting in $W_{0}$. Moreover, such solutions are continuous with respect to the initial data in $B_{Y_{0}}(0, r)$.

## Global Existence

Proof of Theorem 11: By Corollary 19, the problem (1)-(3) has a local solution $\left(u(t), u_{t}(t), v(t), v_{t}(t)\right)$ in $Y_{0}$ defined on some interval $\left[\tau, \tau+t_{0}\right]$.
We can show that

$$
\begin{equation*}
\frac{d}{d t} \mathcal{E}(t)=-\eta\left\|(-\Delta)^{\frac{1}{4}} u_{t}\right\|_{X}^{2}-\eta\left\|(-\Delta)^{\frac{1}{4}} v_{t}\right\|_{X}^{2} \tag{20}
\end{equation*}
$$

for all $\tau<t \leq \tau+t_{0}$, where

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\end{equation*}
$$

for all $\tau<t \leq \tau+t_{0}$, where

$$
\begin{align*}
\mathcal{E}(t) & =\frac{1}{2}\|u(t)\|_{X^{\frac{1}{2}}}^{2}+\frac{1}{2}\|u(t)\|_{X}^{2}+\frac{1}{2}\left\|u_{t}(t)\right\|_{X}^{2} \\
& +\frac{1}{2}\|v(t)\|_{X^{\frac{1}{2}}}^{2}+\frac{1}{2}\left\|v_{t}(t)\right\|_{X}^{2}-\int_{\Omega} \int_{0}^{u} f(s) d s d x \tag{21}
\end{align*}
$$

is the total energy associated with the solution $\left(u(t), u_{t}(t), v(t), v_{t}(t)\right)$ of the problem (1)-(3) in $Y_{0}$.

The identity (20) means that the map $t \mapsto \mathcal{E}(t)$ is monotone decreasing along solutions. Moreover, using the property $\mathcal{E}(t) \leq$ $\mathcal{E}(\tau)$ for all $\tau \leq t \leq \tau+t_{0}$, we can obtain a priori estimate of the solution $\left(u(t), u_{t}(t), v(t), v_{t}(t)\right)$ in $Y_{0}$. In fact, we obtain

$$
\|u\|_{X^{\frac{1}{2}}}^{2}+\left\|u_{t}\right\|_{X}^{2}+\|v\|_{X^{\frac{1}{2}}}^{2}+\left\|v_{t}\right\|_{X}^{2} \leq 4\left(\mathcal{E}(\tau)+C_{\frac{\lambda_{1}}{4}}\right)
$$

that is,

$$
\left\|\left(u(t), u_{t}(t), v(t), v_{t}(t)\right)\right\|_{Y_{0}}^{2} \leq 4\left(\mathcal{E}(\tau)+C_{\frac{\lambda_{1}}{4}}\right)
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$$

that is,

$$
\left\|\left(u(t), u_{t}(t), v(t), v_{t}(t)\right)\right\|_{Y_{0}}^{2} \leq 4\left(\mathcal{E}(\tau)+C_{\frac{\lambda_{1}}{4}}\right)
$$

This ensures that the problem $(1)-(3)$ has a global solution $W(t)$ in $Y_{0}$, which proves the result.

Since the problem (1) - (3) has a global solution $W(t)$ in $Y_{0}$, we can define an evolution process $\{S(t, \tau): t \geq \tau \in \mathbb{R}\}$ in $Y_{0}$ by

$$
\begin{equation*}
S(t, \tau) W_{0}=W(t), \quad t \geq \tau \in \mathbb{R} \tag{22}
\end{equation*}
$$

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S(t, \tau) W_{0}=W(t), \quad t \geq \tau \in \mathbb{R} \tag{22}
\end{equation*}
$$

According to Carvalho and Nascimento (DCDS-S 2009)

$$
\begin{equation*}
S(t, \tau) W_{0}=L(t, \tau) W_{0}+U(t, \tau) W_{0}, \quad t \geq \tau \in \mathbb{R} \tag{23}
\end{equation*}
$$

where $\{L(t, \tau): t \geq \tau \in \mathbb{R}\}$ is the linear evolution process in $Y_{0}$ associated with the homogeneous problem

$$
\left\{\begin{array}{l}
W_{t}+\mathcal{A}(t) W=0, t>\tau  \tag{24}\\
W(\tau)=W_{0}, \tau \in \mathbb{R}
\end{array}\right.
$$

and

$$
\begin{equation*}
U(t, \tau) W_{0}=\int_{\tau}^{t} L(t, s) F\left(S(s, \tau) W_{0}\right) d s \tag{25}
\end{equation*}
$$

## Dissipativeness of the thermoelastic equation

In this section, we prove the existence of the pullback attractor of the problem (1)-(3). To this end, we need to make a modification on the energy functional. More precisely, for $\gamma_{1}, \gamma_{2} \in \mathbb{R}_{+}$, let us define $L_{\gamma_{1}, \gamma_{2}}: Y_{0} \rightarrow \mathbb{R}$ by the map

$$
\begin{align*}
L_{\gamma_{1}, \gamma_{2}}(\phi, \varphi, \psi, \Phi) & =\frac{1}{2}\|\phi\|_{X^{\frac{1}{2}}}^{2}+\frac{1}{2}\|\phi\|_{X}^{2}+\frac{1}{2}\|\varphi\|_{X}^{2}+\frac{1}{2}\|\psi\|_{X^{\frac{1}{2}}}^{2} \\
& +\frac{1}{2}\|\Phi\|_{X}^{2}+\gamma_{1}\langle\phi, \varphi\rangle_{X}+\gamma_{2}\langle\psi, \Phi\rangle_{X} \\
& -\int_{\Omega} \int_{0}^{\phi} f(s) d s d x . \tag{26}
\end{align*}
$$

## Theorem 20

There exists $R>0$ such that for any bounded subset $B \subset Y_{0}$ one can find $t_{0}(B)>0$ satisfying

$$
\left\|\left(u, u_{t}, v, v_{t}\right)\right\|_{Y_{0}}^{2} \leq R \quad \text { for all } \quad t \geq \tau+t_{0}(B)
$$

In particular, the evolution process $\{S(t, \tau): t \geq \tau \in \mathbb{R}\}$ defined in (22) is pullback strongly bounded dissipative.

## Theorem 20

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$$

In particular, the evolution process $\{S(t, \tau): t \geq \tau \in \mathbb{R}\}$ defined in (22) is pullback strongly bounded dissipative.

Next, we prove that the solutions of problem (14) are uniformly exponentially dominated when the initial data are in bounded subsets of $Y_{0}$.

## Theorem 21

Let $B \subset Y_{0}$ be a bounded set. If $W:[\tau, \infty) \rightarrow Y_{0}$ is the global solution of (14) starting at $W_{0} \in B$, then there are positive constants $\sigma=\sigma(B), K_{1}=K_{1}(B)$ and $K_{2}=K_{2}(B)$ such that

$$
\|W(t)\|_{Y_{0}}^{2} \leq K_{1} e^{-\sigma(t-\tau)}+K_{2}, \quad t \geq \tau
$$

## Theorem 22

Let $B \subset Y_{0}$ be a bounded set and denote by $L:[\tau, \infty) \rightarrow Y_{0}$ the solution of the homogeneous problem (24) starting in $W_{0} \in B$. Then there exist positive constants $K=K(B)$ and $\zeta$ such that

$$
\|L(t)\|_{Y_{0}}^{2} \leq K e^{-\zeta(t-\tau)}, \quad t \geq \tau
$$

## Theorem 22

Let $B \subset Y_{0}$ be a bounded set and denote by $L:[\tau, \infty) \rightarrow Y_{0}$ the solution of the homogeneous problem (24) starting in $W_{0} \in B$. Then there exist positive constants $K=K(B)$ and $\zeta$ such that

$$
\|L(t)\|_{Y_{0}}^{2} \leq K e^{-\zeta(t-\tau)}, \quad t \geq \tau
$$

## Proposition 23

For each $t>\tau \in \mathbb{R}$, the evolution process $S(t, \tau): Y_{0} \rightarrow Y_{0}$ given in (22) is a compact map.

Proof of Theorem 12: Theorem 20 assures that the evolution process $S(t, \tau): Y_{0} \rightarrow Y_{0}$ given by (22) is pullback strongly bounded dissipative. Additionally, it follows by Proposition 23 that $S(t, \tau): Y_{0} \rightarrow Y_{0}$ is compact, and, consequently, it is pullback asymptotically compact. Now the result is a simple consequence of Theorem 9.

## Obrigado.

