# EXPONENTIAL ATTRACTORS FOR ABSTRACT EQUATIONS WITH MEMORY AND APPLICATIONS TO VISCOELASTICITY 

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Abstract. We consider an abstract equation with memory of the form

$$
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} k(s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s+\boldsymbol{B} \boldsymbol{x}(t)=0
$$

where $\boldsymbol{A}, \boldsymbol{B}$ are operators acting on some Banach space, and the convolution kernel $k$ is a nonnegative convex summable function of unit mass. The system is translated into an ordinary differential equation on a Banach space accounting for the presence of memory, both in the so-called history space framework and in the minimal state one. The main theoretical result is a theorem providing sufficient conditions in order for the related solution semigroups to possess finite-dimensional exponential attractors. As an application, we prove the existence of exponential attractors for the integrodifferential equation

$$
\partial_{t t} u-h(0) \Delta u-\int_{0}^{\infty} h^{\prime}(s) \Delta u(t-s) \mathrm{d} s+f(u)=g
$$

arising in the theory of isothermal viscoelasticity, which is just a particular concrete realization of the abstract model, having defined the new kernel $h(s)=$ $k(s)+1$.

1. Introduction. A large class of physical phenomena in which delay effects occur, such as viscoelasticity, population dynamics or heat flow in real conductors, are modeled by equations with memory, where the dynamics is influenced by the past history of the variables in play through a convolution integral. Given a real Banach space $X$, the general structure of an equation with memory in the unknown $\boldsymbol{x}=$ $\boldsymbol{x}(t): \mathbb{R} \rightarrow X$ reads as follows:

$$
\begin{equation*}
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} k(s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s+\boldsymbol{B} \boldsymbol{x}(t)=0 \tag{1}
\end{equation*}
$$

[^0]Here, the convolution (or memory) kernel $k$ is a nonnegative summable function of total mass

$$
\int_{0}^{\infty} k(s) \mathrm{d} s=1
$$

having the explicit form

$$
k(s)=\int_{s}^{\infty} \mu(y) \mathrm{d} y
$$

where $\mu \in L^{1}\left(\mathbb{R}^{+}\right)$is a nonincreasing (hence nonnegative) piecewise absolutely continuous function, possibly unbounded about zero. The discontinuity points of $\mu$, if any, form an increasing sequence $s_{n}$, which can be either finite or $s_{n} \rightarrow \infty$. Assuming without loss of generality $\mu$ right-continuous, we denote by

$$
\mu_{n}=\mu\left(s_{n}^{-}\right)-\mu\left(s_{n}\right)>0
$$

the jump amplitudes at the (left) discontinuity points $s_{n}$. Besides, $\mu$ is supposed to satisfy for every $t, s>0$ and some $\Theta \geq 1$ and $\delta>0$ the inequality

$$
\begin{equation*}
\mu(t+s) \leq \Theta^{-\delta t} \mu(s) \tag{2}
\end{equation*}
$$

Concerning $\boldsymbol{A}$ and $\boldsymbol{B}$, they are (possibly nonlinear) operators densely defined on $X$. For any fixed time $t$, the operator $\boldsymbol{B}$ acts on $\boldsymbol{x}(t)$, the instantaneous value of $\boldsymbol{x}$, while $\boldsymbol{A}$ acts on the past history of $\boldsymbol{x}$, namely, the past values of $\boldsymbol{x}$ up to $t$. The variable $\boldsymbol{x}(t)$ is supposed to solve the equation in some weak sense for every $t>0$, whereas it is regarded as a known initial datum for negative times. Accordingly, (1) is supplemented with the "initial condition"

$$
\begin{equation*}
\boldsymbol{x}(-s)=\boldsymbol{h}(s), \quad \forall s \geq 0 \tag{3}
\end{equation*}
$$

where $\boldsymbol{h}:[0, \infty) \rightarrow X$ is a given function accounting for the initial past history of $\boldsymbol{x}$.

A common feature of equations arising from concrete physical models is the presence of dissipation mechanisms. In this perspective, here we are mostly interested in the longterm behavior of solutions. As we will show in a while, it is possible to translate the original problem within a semigroup framework. Accordingly, the dissipativity properties of (1) can be described in terms of "small" sets of the phase space able to eventually capture the trajectories of the related solution semigroup $S(t)$ acting on a suitable Banach space $\mathcal{H}$. Dealing with semigroups, an important object is the global attractor (see e.g. $[1,3,22,23,29,30]$ ), whose existence has been proved for several models with memory. Instead, our main goal is to discuss the existence of exponential attractors, otherwise called inertial sets, firstly introduced in [13] in a Hilbert space setting (see [14] for the Banach space case), which have the advantage of being more stable than global attractors, and attract trajectories exponentially fast (see $[15,24]$ for a detailed discussion). By definition, an exponential attractor for a semigroup $S(t)$ acting on a Banach space $\mathcal{H}$ is a compact set $\mathfrak{E} \subset \mathcal{H}$ satisfying the following properties:

- $\mathfrak{E}$ is positively invariant for the semigroup, namely, $S(t) \mathfrak{E} \subset \mathfrak{E}$ for all $t \geq 0$.
- E has finite fractal dimension in $\mathcal{H}$.
- $\mathfrak{E}$ is exponentially attracting for $S(t)$, i.e. there exist an exponential rate $\omega>0$ and a nondecreasing positive function $Q$ such that

$$
\operatorname{dist}_{\mathcal{H}}(S(t) \mathcal{B}, \mathfrak{E}) \leq Q\left(\|\mathcal{B}\|_{\mathcal{H}}\right) \mathrm{e}^{-\omega t}
$$

for every bounded subset $\mathcal{B} \subset \mathcal{H}$.

Recall that the fractal dimension of $\mathfrak{E}$ in $\mathcal{H}$ is defined as

$$
\operatorname{dim}_{\mathcal{H}}(\mathfrak{E})=\limsup _{r \rightarrow 0} \frac{\ln N(r)}{\ln \frac{1}{r}}
$$

where $N(r)$ is the smallest number of $r$-balls of $\mathcal{H}$ necessary to cover $\mathfrak{E}$, whereas

$$
\operatorname{dist}_{\mathcal{H}}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)=\sup _{b_{1} \in \mathcal{B}_{1}} \inf _{b_{2} \in \mathcal{B}_{2}}\left\|b_{1}-b_{2}\right\|_{\mathcal{H}}
$$

denotes the standard Hausdorff semidistance in $\mathcal{H}$ between two sets $\mathcal{B}_{1}$ and $\mathcal{B}_{2}$.
Indeed, although the existence of global attractors for equations with memory has been investigated in several papers, there are considerably fewer results concerning exponential attractors. This is mainly due to the technical difficulties arising in the application of the classical techniques to this particular framework.

Outline of the paper. As a first step, we show how to translate the original integrodifferential problem (1) into a dynamical system within two different frameworks, the so-called past history framework and the minimal state one. This is done in a heuristic way in the next Sec. 2, and it is then formalized in the proper mathematical setting in Sec. 4, upon defining suitable functional spaces (see Sec. 3 ). This leads to the generation of two strongly continuous semigroups $S(t)$ (on the past history space) and $\hat{S}(t)$ (on the minimal state space) describing the solutions to (1). Sec. 5 is devoted to the main abstract result of this work: namely, we provide sufficient conditions in order for the semigroup $S(t)$ to possess a regular exponential attractor. As a corollary, in Sec. 6 the analogous result is shown to hold for the semigroup $\hat{S}(t)$. In the second part of the paper, we discuss an application of the abstract theorems to the well-known equation of viscoelasticity (see Sec. 7 and Sec. $8)$.
2. The transformed equation. In order to view (1) as a dynamical system, two different strategies have been devised.
2.1. The past history framework. Using a method proposed by Dafermos [9] (see also [18]), we introduce for $(t, s) \in[0, \infty) \times \mathbb{R}^{+}$the summed past history

$$
\boldsymbol{\eta}^{t}(s)=\int_{0}^{s} \boldsymbol{A} \boldsymbol{x}(t-y) \mathrm{d} y
$$

which (formally) fulfills the differential equation of hyperbolic type

$$
\partial_{t} \boldsymbol{\eta}^{t}(s)=-\partial_{s} \boldsymbol{\eta}^{t}(s)+\boldsymbol{A} \boldsymbol{x}(t)
$$

subject to the "boundary condition"

$$
\lim _{s \rightarrow 0} \boldsymbol{\eta}^{t}(s)=0
$$

A formal integration by parts yields the equality

$$
\int_{0}^{\infty} k(s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s=\int_{0}^{\infty} \mu(s) \boldsymbol{\eta}^{t}(s) \mathrm{d} s
$$

Hence, the original equation (1) translates into the evolution system in the unknown variables $\boldsymbol{x}=\boldsymbol{x}(t)$ and $\boldsymbol{\eta}=\boldsymbol{\eta}^{t}(s)$

$$
\left\{\begin{array}{l}
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} \mu(s) \boldsymbol{\eta}^{t}(s) \mathrm{d} s+\boldsymbol{B} \boldsymbol{x}(t)=0  \tag{4}\\
\partial_{t} \boldsymbol{\eta}^{t}(s)=-\partial_{s} \boldsymbol{\eta}^{t}(s)+\boldsymbol{A} \boldsymbol{x}(t)
\end{array}\right.
$$

In turn, the initial condition (3) becomes

$$
\left\{\begin{array}{l}
\boldsymbol{x}(0)=\boldsymbol{x}_{0} \\
\boldsymbol{\eta}^{0}=\boldsymbol{\eta}_{0}
\end{array}\right.
$$

where we set

$$
\boldsymbol{x}_{0}=\boldsymbol{h}(0) \quad \text { and } \quad \boldsymbol{\eta}_{0}(s)=\int_{0}^{s} \boldsymbol{A} \boldsymbol{h}(y) \mathrm{d} y
$$

Remark 1. If $\boldsymbol{A}$ is linear, and so it commutes with the integral, one usually defines $\boldsymbol{\eta}$ in a slightly different way, namely, $\boldsymbol{\eta}^{t}(s)=\int_{0}^{s} \boldsymbol{x}(t-y) \mathrm{d} y$. Clearly, (4) changes accordingly.
2.2. The minimal state framework. Although it successfully allows to view the original integrodifferential equation as an abstract differential equation, the past history framework suffers from a structural theoretical drawback. Indeed, it is clear that the dynamics is known once $\boldsymbol{x}_{0}$ and $\boldsymbol{\eta}_{0}$ are assigned. This means that the initial past history $\boldsymbol{h}$ may not be recovered from the future evolutions, hence it has no possibilities to be an observable quantity. To cope with this fact, an alternative approach has been proposed (see [10, 11, 17]), based on the notion of minimal state: an additional variable accounting for the past history which contains the necessary and sufficient information determining the future dynamics. Precisely, for $(t, \tau) \in[0, \infty) \times \mathbb{R}^{+}$, we define

$$
\boldsymbol{\xi}^{t}(\tau)=\int_{0}^{\infty} \mu(\tau+s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s
$$

which (again, formally) satisfies the relations

$$
\partial_{t} \boldsymbol{\xi}^{t}(\tau)=\partial_{\tau} \boldsymbol{\xi}^{t}(\tau)+\mu(\tau) \boldsymbol{A} \boldsymbol{x}(t)
$$

and

$$
\int_{0}^{\infty} \boldsymbol{\xi}^{t}(\tau) \mathrm{d} \tau=\int_{0}^{\infty} k(s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s
$$

Hence, (1) takes the form

$$
\left\{\begin{array}{l}
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} \boldsymbol{\xi}^{t}(\tau) \mathrm{d} \tau+\boldsymbol{B} \boldsymbol{x}(t)=0  \tag{5}\\
\partial_{t} \boldsymbol{\xi}^{t}(\tau)=\partial_{\tau} \boldsymbol{\xi}^{t}(\tau)+\mu(\tau) \boldsymbol{A} \boldsymbol{x}(t)
\end{array}\right.
$$

and the initial condition (3) translates into

$$
\left\{\begin{array}{l}
\boldsymbol{x}(0)=\boldsymbol{x}_{0} \\
\boldsymbol{\xi}^{0}=\boldsymbol{\xi}_{0}
\end{array}\right.
$$

where we set

$$
\boldsymbol{x}_{0}=\boldsymbol{h}(0) \quad \text { and } \quad \boldsymbol{\xi}_{0}(\tau)=\int_{0}^{\infty} \mu(\tau+s) \boldsymbol{A} \boldsymbol{h}(s) \mathrm{d} s
$$

Remark 2. Such a description meets the sought minimality requirement. Indeed, once the initial state $\boldsymbol{\xi}_{0}$ is assigned, we can write

$$
\boldsymbol{\xi}^{t}(\tau)=\boldsymbol{\xi}_{0}(t+\tau)+\int_{0}^{t} \mu(\tau+s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s
$$

Then, plugging the latter equality into the first equation of (5), we deduce for every $t \geq 0$ the relation

$$
\int_{t}^{\infty} \boldsymbol{\xi}_{0}(\tau) \mathrm{d} \tau=G(t)
$$

for some function $G$ depending only on the values of $\boldsymbol{x}(t)$ for $t$ positive. Hence, the knowledge of $\boldsymbol{x}(t)$ for all $t \geq 0$ uniquely determines $\boldsymbol{\xi}_{0}$.
Remark 3. Similarly to the past history case, if $\boldsymbol{A}$ is linear it is customary to define instead $\boldsymbol{\xi}^{t}(\tau)=\int_{0}^{\infty} \mu(\tau+s) \boldsymbol{x}(t-s) \mathrm{d} s$, and change (5) accordingly.
3. Functional setting and notation. We first introduce the abstract functional setting needed to carry out our analysis.
3.1. Notation. We denote by $\mathfrak{I}$ the set of nondecreasing functions $Q:[0, \infty) \rightarrow$ $[0, \infty)$, and by $\mathfrak{D}$ the set of nonincreasing functions $q:[0, \infty) \rightarrow[0, \infty)$ vanishing at infinity. Besides, given a Banach space $\mathcal{H}$ and $r \geq 0$, we set

$$
\mathbb{B}_{\mathcal{H}}(r)=\left\{z \in \mathcal{H}:\|z\|_{\mathcal{H}} \leq r\right\}
$$

3.2. Geometric spaces. Let $X^{0}, X^{1}$ and $Y^{0}, Y^{1}$ be reflexive Banach spaces having dense and compact embeddings

$$
X^{1} \Subset X^{0} \quad \text { and } \quad Y^{1} \Subset Y^{0}
$$

3.3. Memory spaces. For $\imath=0,1$, we introduce the memory spaces

$$
\mathcal{M}^{\imath}=L_{\mu}^{2}\left(\mathbb{R}^{+} ; Y^{\imath}\right)
$$

namely, the spaces of $L^{2}$-functions on $\mathbb{R}^{+}$with values in $Y^{\imath}$ with respect to the measure $\mu(s) \mathrm{d} s$, normed by

$$
\|\boldsymbol{\eta}\|_{\mathcal{M}^{2}}^{2}=\int_{0}^{\infty} \mu(s)\|\boldsymbol{\eta}(s)\|_{Y^{2}}^{2} \mathrm{~d} s
$$

If $Y^{\imath}$ is a Hilbert space, so is $\mathcal{M}^{2}$. Albeit clearly continuous, the embedding $\mathcal{M}^{1} \subset$ $\mathcal{M}^{0}$ is not compact (see [27] for a counterexample). We will also consider the linear operator

$$
T \boldsymbol{\eta}=-\boldsymbol{\eta}^{\prime}, \quad \operatorname{dom}(T)=\left\{\boldsymbol{\eta} \in \mathcal{M}^{0}: \boldsymbol{\eta}^{\prime} \in \mathcal{M}^{0}, \boldsymbol{\eta}(0)=0\right\}
$$

where the prime stands for weak derivative and $\boldsymbol{\eta}(0)=\lim _{s \rightarrow 0} \boldsymbol{\eta}(s)$ in $Y^{0}$. It can be verified that $T$ is the infinitesimal generator of the right-translation semigroup $R(t)$ on $\mathcal{M}^{0}$, defined as

$$
(R(t) \boldsymbol{\eta})(s)= \begin{cases}0 & s \leq t \\ \boldsymbol{\eta}(s-t) & s>t\end{cases}
$$

Finally, for $\imath=0,1$, we define the extended memory spaces

$$
\mathcal{H}^{\imath}=X^{\imath} \times \mathcal{M}^{\imath}
$$

which are Banach spaces with respect to the norms

$$
\|(\boldsymbol{x}, \boldsymbol{\eta})\|_{\mathcal{H}^{2}}^{2}=\|\boldsymbol{x}\|_{X^{2}}^{2}+\|\boldsymbol{\eta}\|_{\mathcal{M}^{2}}^{2}
$$

and we call $\Pi_{1}$ and $\Pi_{2}$ the projections of $\mathcal{H}^{0}$ onto its components $X^{0}$ and $\mathcal{M}^{0}$, namely,

$$
\Pi_{1}(\boldsymbol{x}, \boldsymbol{\eta})=\boldsymbol{x}, \quad \Pi_{2}(\boldsymbol{x}, \boldsymbol{\eta})=\boldsymbol{\eta}
$$

3.4. State spaces. We define the nonnegative kernel $\nu$ as

$$
\nu(\tau)=1 / \mu(\tau)
$$

where $\nu(0)=\lim _{\tau \rightarrow 0} 1 / \mu(\tau)$, and we agree to set $\nu(\tau)=0$ whenever $\mu(\tau)=0$, in order to include the finite delay case in our discussion. The assumptions on $\mu$ imply that $\nu$ is nondecreasing and piecewise absolutely continuous. In particular, (2) yields

$$
\begin{equation*}
\nu(\tau-s) \leq \Theta \mathrm{e}^{-\delta s} \nu(\tau), \quad \forall s<\tau \tag{6}
\end{equation*}
$$

provided that $\nu(\tau)>0$. Then, for $\imath=0,1$, we introduce the (minimal) state spaces

$$
\mathcal{S}^{\imath}=L_{\nu}^{2}\left(\mathbb{R}^{+} ; Y^{\imath}\right)
$$

with norms

$$
\|\boldsymbol{\xi}\|_{\mathcal{S}^{2}}^{2}=\int_{0}^{\infty} \nu(\tau)\|\boldsymbol{\xi}(\tau)\|_{Y^{2}}^{2} \mathrm{~d} \tau
$$

along with the extended state spaces

$$
\mathcal{V}^{\imath}=X^{\imath} \times \mathcal{S}^{\imath}
$$

normed by

$$
\|(\boldsymbol{x}, \boldsymbol{\xi})\|_{\mathcal{V}^{2}}^{2}=\|\boldsymbol{x}\|_{X^{2}}^{2}+\|\boldsymbol{\xi}\|_{\mathcal{S}^{2}}^{2} .
$$

Besides, let $P$ be the infinitesimal generator of the left-translation semigroup on $\mathcal{S}^{0}$, namely,

$$
P \boldsymbol{\xi}=\boldsymbol{\xi}^{\prime}, \quad \operatorname{dom}(P)=\left\{\boldsymbol{\xi} \in \mathcal{S}^{0}: \boldsymbol{\xi}^{\prime} \in \mathcal{S}^{0}\right\}
$$

Abusing the notation, we keep denoting by $\Pi_{1}$ and $\Pi_{2}$ the projections of $\mathcal{V}^{0}$ onto its components $X^{0}$ and $\mathcal{S}^{0}$, respectively.
4. The solution semigroups. We are now in the position to rephrase (4) and (5) within the correct functional setting.
4.1. The past history framework. We interpret (4) as the ODE in $\mathcal{H}^{0}$ in the unknown variables $\boldsymbol{x}=\boldsymbol{x}(t)$ and $\boldsymbol{\eta}=\boldsymbol{\eta}^{t}(s)$

$$
\left\{\begin{array}{l}
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} \mu(s) \boldsymbol{\eta}^{t}(s) \mathrm{d} s+\boldsymbol{B} \boldsymbol{x}(t)=0  \tag{7}\\
\partial_{t} \boldsymbol{\eta}^{t}=T \boldsymbol{\eta}^{t}+\boldsymbol{A} \boldsymbol{x}(t)
\end{array}\right.
$$

Throughout this work, we assume that (7) has a unique global solution in some weak sense for every initial datum $\boldsymbol{z}=\left(\boldsymbol{x}_{0}, \boldsymbol{\eta}_{0}\right) \in \mathcal{H}^{0}$. This is the same as saying that (7) generates a semigroup

$$
S(t): \mathcal{H}^{0} \rightarrow \mathcal{H}^{0}
$$

We do not require any continuity in time, although in most concrete cases one has

$$
t \mapsto S(t) \boldsymbol{z} \in \mathcal{C}\left([0, \infty), \mathcal{H}^{0}\right), \quad \forall \boldsymbol{z} \in \mathcal{H}^{0}
$$

In particular, once $\boldsymbol{x}(t)$ is known, the second component $\boldsymbol{\eta}^{t}=\Pi_{2} S(t) \boldsymbol{z}$ of the solution is recovered as a Duhamel integral, hence it fulfills the explicit representation formula

$$
\begin{equation*}
\boldsymbol{\eta}^{t}(s)=\left(R(t) \boldsymbol{\eta}_{0}\right)(s)+\int_{0}^{\min \{t, s\}} \boldsymbol{A} \boldsymbol{x}(t-y) \mathrm{d} y \tag{8}
\end{equation*}
$$

Besides, $S(t)$ is supposed to satisfy the following Hölder continuity property: there exists $Q \in \mathfrak{I}$ and $\kappa=\kappa(r) \in(0,1]$ such that

$$
\begin{equation*}
\left\|S(t) \boldsymbol{z}_{1}-S(t) \boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}} \leq Q(r) \mathrm{e}^{Q(r) t}\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{\kappa} \tag{9}
\end{equation*}
$$

whenever $\boldsymbol{z}_{1}, \boldsymbol{z}_{2} \in \mathbb{B}_{\mathcal{H}^{0}}(r)$.
4.2. The minimal state framework. In a similar manner, we interpret (5) as the ODE in $\mathcal{V}^{0}$ in the unknown variables $\boldsymbol{x}=\boldsymbol{x}(t)$ and $\boldsymbol{\xi}=\boldsymbol{\xi}^{t}(\tau)$

$$
\left\{\begin{array}{l}
\partial_{t} \boldsymbol{x}(t)+\int_{0}^{\infty} \boldsymbol{\xi}^{t}(\tau) \mathrm{d} \tau+\boldsymbol{B} \boldsymbol{x}(t)=0  \tag{10}\\
\partial_{t} \boldsymbol{\xi}^{t}=P \boldsymbol{\xi}^{t}+\mu \boldsymbol{A} \boldsymbol{x}(t)
\end{array}\right.
$$

Again, for every initial datum $\hat{\boldsymbol{z}}=\left(\boldsymbol{x}_{0}, \boldsymbol{\xi}_{0}\right) \in \mathcal{V}^{0}$, we assume that (10) generates a semigroup

$$
\hat{S}(t): \mathcal{V}^{0} \rightarrow \mathcal{V}^{0}
$$

Hence, arguing as in the previous case, the second component $\boldsymbol{\xi}^{t}=\Pi_{2} \hat{S}(t) \hat{\boldsymbol{z}}$ of the solution fulfills the explicit representation formula

$$
\begin{equation*}
\boldsymbol{\xi}^{t}(\tau)=\boldsymbol{\xi}_{0}(t+\tau)+\int_{0}^{t} \mu(\tau+s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s \tag{11}
\end{equation*}
$$

Finally, we suppose that there exists $Q \in \mathfrak{I}$ and $\kappa=\kappa(r) \in(0,1]$ such that

$$
\begin{equation*}
\left\|\hat{S}(t) \hat{\boldsymbol{z}}_{1}-\hat{S}(t) \hat{\boldsymbol{z}}_{2}\right\|_{\mathcal{V}^{0}} \leq Q(r) \mathrm{e}^{Q(r) t}\left\|\hat{\boldsymbol{z}}_{1}-\hat{\boldsymbol{z}}_{2}\right\|_{\mathcal{V}^{0}}^{\kappa} \tag{12}
\end{equation*}
$$

whenever $\hat{\boldsymbol{z}}_{1}, \hat{\boldsymbol{z}}_{2} \in \mathbb{B}_{\mathcal{V}^{0}}(r)$.
4.3. The $\operatorname{map} \boldsymbol{\Lambda}$. The connection between the memory and the state spaces has been devised in [5, 17]. Setting for any $\boldsymbol{\eta} \in \mathcal{M}^{0}$

$$
(\Lambda \boldsymbol{\eta})(\tau)=-\int_{0}^{\infty} \mu^{\prime}(\tau+s) \boldsymbol{\eta}(s) \mathrm{d} s+\sum_{\tau<s_{n}} \mu_{n} \boldsymbol{\eta}\left(s_{n}-\tau\right)
$$

the following lemma is proved.
Lemma 4.1. The map $\boldsymbol{\Lambda}$ defined by

$$
(\boldsymbol{x}, \boldsymbol{\eta}) \mapsto \boldsymbol{\Lambda}(\boldsymbol{x}, \boldsymbol{\eta})=(\boldsymbol{x}, \Lambda \boldsymbol{\eta})
$$

is a bounded linear operator of unitary norm from $\mathcal{H}^{\imath}$ into $\mathcal{V}^{2}$. Moreover, for every $\boldsymbol{\eta} \in \mathcal{M}^{0}$ and every $\tau>0$, we have the equality

$$
\begin{equation*}
\int_{0}^{\infty} \mu(\tau+s) \boldsymbol{\eta}(s) \mathrm{d} s=\int_{\tau}^{\infty}(\Lambda \boldsymbol{\eta})(y) \mathrm{d} y=(\Lambda \boldsymbol{H})(\tau) \tag{13}
\end{equation*}
$$

where $\boldsymbol{H}(s)=\int_{0}^{s} \boldsymbol{\eta}(y) \mathrm{d} y$.
Remark 4. However, it should be observed that, in general, the map $\boldsymbol{\Lambda}$ is neither 1-to-1 nor onto.

The link between the two formulations is detailed in the next lemma (cf. [5, 17]), showing in particular that the state approach describes the dynamics in greater generality.

Lemma 4.2. For every $\boldsymbol{z} \in \mathcal{H}^{0}$, the following equality holds:

$$
\hat{S}(t) \boldsymbol{\Lambda} \boldsymbol{z}=\boldsymbol{\Lambda} S(t) \boldsymbol{z}
$$

## 5. Exponential attractors: The past history framework.

5.1. Statement of the theorem. Our main goal is to provide a "user friendly" recipe establishing sufficient conditions in order for the semigroup $S(t)$ acting on $\mathcal{H}^{0}$ to possess a regular exponential attractor $\mathfrak{E}$. Such a result is specifically tailored for equations with memory: as it will be shown in the final application, the abstract hypotheses are easily translated in concrete terms. A similar attempt has been made in [19], but only for a particular case of our general equation (1). In principle, this theorem can be inferred from the previous work [21], dealing with the existence of a Hölder continuous family of exponential attractors $\mathfrak{E}_{\varepsilon}$ for a family of semigroups $S_{\varepsilon}(t)$ depending on a perturbation parameter $\varepsilon \in[0,1]$; see [21, Theorem 4.4]. Nonetheless, for the basic (but still very interesting) case of a single semigroup of memory type, it is not at all immediate how to derive a clean and simple statement of wide application from [21]. Accordingly, we will give the proof in some detail.

Theorem 5.1. In addition to our general assumptions, suppose that the following hold:
(i) There exists $R_{1}>0$ such that, given any $r \geq 0$,

$$
\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{1}} \leq q_{r}(t)+R_{1}
$$

for some $q_{r} \in \mathfrak{D}$ and every $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{1}}(r)$.
(ii) There is $r_{1}>0$ such that the ball $\mathbb{B}_{\mathcal{H}^{1}}\left(r_{1}\right)$ is exponentially attracting for $S(t)$.
(iii) The operator $\boldsymbol{A}$ maps bounded subsets of $X^{1}$ into bounded subsets of $Y^{0}$.
(iv) There exist $p \in(1, \infty]$ and $Q \in \mathfrak{I}$ such that, for every $r \geq 0$ and every $\theta>0$ sufficiently large,

$$
\left\|\partial_{t} \boldsymbol{x}\right\|_{L^{p}\left(\theta, 2 \theta ; X^{0}\right)} \leq Q(r+\theta)
$$

for all $\boldsymbol{x}(t)=\Pi_{1} S(t) \boldsymbol{z}$ with $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{1}}(r)$.
$(\mathbf{v})$ For every $r>0$ there are $Q_{r} \in \mathfrak{I}$ and $q_{r} \in \mathfrak{D}$ such that, for all $\boldsymbol{z}_{1}, \boldsymbol{z}_{2} \in \mathbb{B}_{\mathcal{H}^{1}}(r)$,

$$
S(t) \boldsymbol{z}_{1}-S(t) \boldsymbol{z}_{2}=L\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)+K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)
$$

where the maps $L$ and $K$ satisfy

$$
\begin{aligned}
\left\|L\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{0}} & \leq q_{r}(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}} \\
\left\|K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{1}} & \leq Q_{r}(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}
\end{aligned}
$$

Moreover, let $\overline{\boldsymbol{\eta}}^{t}=\Pi_{2} K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)$ fulfill the Cauchy problem

$$
\left\{\begin{array}{l}
\partial_{t} \overline{\boldsymbol{\eta}}^{t}=T \overline{\boldsymbol{\eta}}^{t}+\boldsymbol{w}(t) \\
\overline{\boldsymbol{\eta}}^{0}=0
\end{array}\right.
$$

for some function $\boldsymbol{w}$ satisfying the estimate

$$
\|\boldsymbol{w}(t)\|_{Y^{0}} \leq Q_{r}(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}
$$

Then $S(t)$ has an exponential attractor $\mathfrak{E}$ contained in $\mathbb{B}_{\mathcal{H}^{1}}\left(r_{1}\right)$.
5.2. Proof of Theorem 5.1. A first difficulty, in connection with the higher order estimate of point (v), comes from the lack of compactness of the embedding $\mathcal{H}^{1} \subset$ $\mathcal{H}^{0}$, for the embedding $\mathcal{M}^{1} \subset \mathcal{M}^{0}$ is not compact either. In order to bypass this obstacle, we introduce the more regular memory space

$$
\mathcal{K}=\left\{\boldsymbol{\eta} \in \mathcal{M}^{1}: \boldsymbol{\eta}^{\prime} \in \mathcal{M}^{0}, \boldsymbol{\eta}(0)=0, \sup _{y \geq 1} y \mathbb{T}(y ; \boldsymbol{\eta})<\infty\right\}
$$

where

$$
\mathbb{T}(y ; \boldsymbol{\eta})=\int_{y}^{\infty} \mu(s)\|\boldsymbol{\eta}(s)\|_{Y^{0}}^{2} \mathrm{~d} s, \quad y \geq 1
$$

denotes the tail function of $\boldsymbol{\eta}$. Setting

$$
\mathbb{H}(\boldsymbol{\eta})=\left\|\boldsymbol{\eta}^{\prime}\right\|_{\mathcal{M}^{0}}^{2}+\sup _{y \geq 1} y \mathbb{T}(y ; \boldsymbol{\eta})
$$

$\mathcal{K}$ turns out to be a Banach space with the norm

$$
\|\boldsymbol{\eta}\|_{\mathcal{K}}^{2}=\|\boldsymbol{\eta}\|_{\mathcal{M}^{1}}^{2}+\mathbb{H}(\boldsymbol{\eta})
$$

so that $\mathcal{K}$ is continuously embedded into $\mathcal{M}^{1}$. Although $\mathcal{K}$ might not be reflexive, it can be proved that its closed balls are closed in $\mathcal{M}^{0}$ (see [21, Proposition 5.6]). Moreover, owing to [27, Lemma 5.5], we have the compact embedding $\mathcal{K} \Subset \mathcal{M}^{0}$. Finally, we introduce the product space

$$
\mathcal{W}=X^{1} \times \mathcal{K} \Subset \mathcal{H}^{0}
$$

We will need a technical lemma, which easily follows (up to inessential modifications) by collecting Lemma 3.3 and Lemma 3.4 in [8].

Lemma 5.2. Given $\mathrm{T} \in(0, \infty]$, let $\boldsymbol{\eta}_{0} \in \mathcal{M}^{0}$ be such that $\boldsymbol{\eta}_{0}(0)=0$ and $\mathbb{H}\left(\boldsymbol{\eta}_{0}\right)<$ $\infty$, and let $\boldsymbol{\eta}=\boldsymbol{\eta}^{t}(s)$ be the solution to the Cauchy problem on $I_{\mathrm{T}}$ (where $I_{\mathrm{T}}=[0, \mathrm{~T}]$ if $\mathrm{T}<\infty$ and $I_{\mathrm{T}}=[0, \infty)$ otherwise)

$$
\left\{\begin{array}{l}
\partial_{t} \boldsymbol{\eta}^{t}=T \boldsymbol{\eta}^{t}+\boldsymbol{w}(t) \\
\boldsymbol{\eta}^{0}=\boldsymbol{\eta}_{0}
\end{array}\right.
$$

Assume that there exist $q \in \mathfrak{D}$ and $K \geq 0$ such that

$$
\|\boldsymbol{w}(t)\|_{Y^{0}}^{2} \leq q(t)+K
$$

Then, for all $t \in I_{\mathrm{T}}$, we have that $\boldsymbol{\eta}^{t}(0)=0$ and

$$
\mathbb{H}\left(\boldsymbol{\eta}^{t}\right) \leq C_{0} \mathrm{e}^{-\frac{\delta}{2} t}+M K
$$

Here, $C_{0} \geq 0$ depends on $\mathbb{H}\left(\boldsymbol{\eta}_{0}\right)$, $q$ and $K$, whereas $M>0$ is a universal constant and $\delta>0$ is given by (2). In particular, if $\boldsymbol{\eta}_{0}=0$ and $q \equiv 0$, then $C_{0}=0$.

We now proceed to the proof by stating two more lemmas.
Lemma 5.3. There exists $R_{2}>0$ such that, given any $r \geq 0$,

$$
\|S(t) \boldsymbol{z}\|_{\mathcal{W}} \leq q_{r}(t)+R_{2}
$$

for some $q_{r} \in \mathfrak{D}$ and every $\boldsymbol{z} \in \mathbb{B}_{\mathcal{W}}(r)$.
Proof. Let $r \geq 0$, and let $\boldsymbol{z}=\left(\boldsymbol{x}_{0}, \boldsymbol{\eta}_{0}\right) \in \mathbb{B}_{\mathcal{W}}(r)$ be arbitrarily chosen. The inclusion $\mathbb{B}_{\mathcal{W}}(r) \subset \mathbb{B}_{\mathcal{H}^{1}}(r)$ together with (i) entail

$$
\|\boldsymbol{x}(t)\|_{X^{1}} \leq q_{r}(t)+R_{1}
$$

for some $q_{r} \in \mathfrak{D}$. Owing to (iii), up to possibly changing $q_{r}$ and $R_{1}$, the inequality

$$
\|\boldsymbol{A} \boldsymbol{x}(t)\|_{Y^{0}} \leq q_{r}(t)+R_{1}
$$

is readily seen to hold. Since $\boldsymbol{\eta}^{t}$ solves the second equation of system (7), exploiting Lemma 5.2 with $\mathrm{T}=\infty$ we obtain

$$
\mathbb{H}\left(\boldsymbol{\eta}^{t}\right) \leq C_{0} \mathrm{e}^{-\frac{\delta}{2} t}+M R_{1},
$$

where here $C_{0}$ depends on $r$. Recalling (i), this yields the desired conclusion.
Lemma 5.4. There is $r_{2}>0$ such that $\mathbb{B}_{\mathcal{W}}\left(r_{2}\right)$ is exponentially attracting for $S(t)$.
Proof. It is enough to show that the ball $\mathbb{B}_{\mathcal{W}}\left(r_{2}\right)$ exponentially attracts $\mathbb{B}_{\mathcal{H}^{1}}\left(r_{1}\right)$. Indeed, exploiting (ii) and the continuity (9), we can apply the transitivity property of exponential attraction [16], and conclude that the basin of (exponential) attraction of $\mathbb{B}_{\mathcal{W}}\left(r_{2}\right)$ coincides with the whole phase space. Let then $\left(\boldsymbol{x}(t), \boldsymbol{\eta}^{t}\right)$ be the solution with initial datum $\boldsymbol{z}=\left(\boldsymbol{x}_{0}, \boldsymbol{\eta}_{0}\right) \in \mathbb{B}_{\mathcal{H}^{1}}\left(r_{1}\right)$. Along the proof, $Q \in \mathfrak{I}$ denotes a generic function. Making use of (i) and (iii), we get

$$
\sup _{t \geq 0}\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{1}} \leq Q\left(r_{1}\right) \quad \Rightarrow \quad \sup _{t \geq 0}\|\boldsymbol{A} \boldsymbol{x}(t)\|_{Y^{0}} \leq Q\left(r_{1}\right)
$$

We make the decomposition

$$
S(t) \boldsymbol{z}=\left(0, \boldsymbol{\psi}^{t}\right)+\left(\boldsymbol{x}(t), \boldsymbol{\zeta}^{t}\right)
$$

where $\boldsymbol{\psi}^{t}$ and $\boldsymbol{\zeta}^{t}$ solve

$$
\left\{\begin{array} { l } 
{ \partial _ { t } \boldsymbol { \psi } ^ { t } = T \boldsymbol { \psi } ^ { t } , } \\
{ \boldsymbol { \psi } ^ { 0 } = \boldsymbol { \eta } _ { 0 } , }
\end{array} \quad \left\{\begin{array}{l}
\partial_{t} \boldsymbol{\zeta}^{t}=T \boldsymbol{\zeta}^{t}+\boldsymbol{A} \boldsymbol{x}(t) \\
\boldsymbol{\zeta}^{0}=0
\end{array}\right.\right.
$$

Condition (2) entails that $\boldsymbol{\psi}^{t}$ satisfies the uniform exponential decay

$$
\left\|\boldsymbol{\psi}^{t}\right\|_{\mathcal{M}^{1}}^{2} \leq \Theta r_{1}^{2} \mathrm{e}^{-\delta t}
$$

In particular, $\boldsymbol{\psi}^{t}$ is uniformly bounded in $\mathcal{M}^{1}$, the bound depending on $r_{1}$. Since so is $\boldsymbol{\eta}^{t}$, we learn that

$$
\sup _{t \geq 0}\left\|\boldsymbol{\zeta}^{t}\right\|_{\mathcal{M}^{1}} \leq Q\left(r_{1}\right)
$$

which, combined with an application of Lemma 5.2 with $\mathrm{T}=\infty$, gives

$$
\left\|\boldsymbol{\zeta}^{t}\right\|_{\mathcal{K}}^{2}=\left\|\boldsymbol{\zeta}^{t}\right\|_{\mathcal{M}^{1}}^{2}+\mathbb{H}\left(\boldsymbol{\zeta}^{t}\right) \leq Q\left(r_{1}\right)
$$

We conclude that

$$
\left\|\left(\boldsymbol{x}(t), \boldsymbol{\zeta}^{t}\right)\right\|_{\mathcal{W}} \leq Q\left(r_{1}\right)
$$

while $\left(0, \boldsymbol{\psi}^{t}\right)$ decays exponentially (in fact, even in $\left.\mathcal{H}^{1}\right)$. Hence, the claim follows by choosing $r_{2}=r_{2}\left(r_{1}\right)$ sufficiently large.

At this point, with reference to Lemmas 5.3 and 5.4 , choose $\varrho>\max \left\{r_{2}, R_{2}\right\}$ and consider the ball

$$
\mathbb{B}=\mathbb{B}_{\mathcal{W}}(\varrho)
$$

Since $\varrho>r_{2}$, it is clear that $\mathbb{B}$ remains exponentially attracting for $S(t)$; whereas, since $\varrho>R_{2}$, we have for $t^{*}>0$ large enough

$$
q_{\varrho}\left(t^{*}\right)+R_{2} \leq \varrho,
$$

implying in turn

$$
S(t) \mathbb{B} \subset \mathbb{B}, \quad \forall t \geq t^{*}
$$

Fixing a positive $\lambda<1$ and exploiting point (v), up to possibly increasing $t^{*}$, the map

$$
S=S\left(t^{*}\right): \mathbb{B} \rightarrow \mathbb{B}
$$

fulfills for every $\boldsymbol{z}_{1}, \boldsymbol{z}_{2} \in \mathbb{B}$ the decomposition

$$
S \boldsymbol{z}_{1}-S \boldsymbol{z}_{2}=L\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)+K\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)
$$

where

$$
\begin{aligned}
& \left\|L\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{0}} \leq \lambda\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}} \\
& \left\|K\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{1}} \leq C\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}
\end{aligned}
$$

Here and in what follows, $C>0$ is a generic constant depending on $\varrho$ and $t^{*}$. Besides, applying Lemma 5.2 with $\mathrm{T}=t^{*}$, we get

$$
\mathbb{H}\left(\overline{\boldsymbol{\eta}}^{t^{*}}\right) \leq C\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2}
$$

Hence,

$$
\left\|K\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{W}} \leq C\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}} .
$$

Invoking a nowadays classical result ${ }^{1}$ (see [4, 14]), there exists an exponential attractor $\mathfrak{E}_{\mathrm{d}} \subset \mathbb{B}$ for the discrete semigroup $S^{n}$ made by the $n^{\text {th }}$-iterations of $S$. Then, we define the map $\mathcal{S}:\left[t^{*}, 2 t^{*}\right] \times \mathbb{B} \rightarrow \mathbb{B}$ by

$$
\mathcal{S}(t, \boldsymbol{z})=S(t) \boldsymbol{z}
$$

and we consider the set

$$
\mathfrak{E}=\mathcal{S}\left(\left[t^{*}, 2 t^{*}\right] \times \mathfrak{E}_{\mathrm{d}}\right) \subset \mathbb{B}
$$

It is apparent that $\mathfrak{E}$ is positively invariant for $S(t)$. Besides, since $\boldsymbol{\eta}^{t}$ is a solution to the second equation of system (7), using (iii) and arguing exactly as in [21, Proposition 7.2], we get

$$
\left\|\boldsymbol{\eta}^{t_{1}}-\boldsymbol{\eta}^{t_{2}}\right\|_{\mathcal{M}^{0}} \leq C\left(t_{1}-t_{2}\right), \quad \forall \boldsymbol{z} \in \mathbb{B}
$$

for every $t_{1}>t_{2} \geq t^{*}$. This estimate together with point (iv) and (9) imply that the map $\mathcal{S}$ is Hölder continuous when $\mathbb{B}$ is endowed with the $\mathcal{H}^{0}$-topology. Therefore, $\mathfrak{E}$ is compact in $\mathcal{H}^{0}$ and, due to the well-known properties of the fractal measure, it has finite fractal dimension in $\mathcal{H}^{0}$. By the continuity property (9), we first deduce that $\mathfrak{E}$ exponentially attracts $\mathbb{B}$ under the action of the continuous semigroup $S(t)$, and appealing once again to the transitivity property of exponential attraction [16], we conclude that $\mathfrak{E}$ is exponentially attracting on the whole space $\mathcal{H}^{0}$.

As a final comment, we point out that one interesting feature of the theorem is that the higher regularity memory space $\mathcal{K}$ does not appear in the statement, being completely hidden in the proof. In particular, in concrete cases one has to work only with the more natural spaces $\mathcal{M}^{2}$. This renders the result more immediate and easier to apply.
5.3. Further remarks. We conclude the section with some remarks on the assumptions of Theorem 5.1.
I. Hypothesis (2) on $\mu$ cannot be weakened since, as proved in [2], it is necessary (and sufficient as well) for the uniform stability of the semigroup $R(t)$ on $\mathcal{M}^{0}$. Indeed, in order to have existence of absorbing sets for systems with memory, postulated by assumptions (i) and (ii), it is necessary that $R(t)$ be uniformly stable. We also mention that, when $\Theta=1$, hypothesis (2) boils down the well-known condition devised by Dafermos [9]

$$
\begin{equation*}
\mu^{\prime}(s)+\delta \mu(s) \leq 0 \tag{14}
\end{equation*}
$$

[^1]On the other hand, when $\Theta>1$, the gap between (2) and the above relation becomes quite effective. Just note that, contrary to (2), the latter condition does not allow $\mu$ to have flat zones, or even horizontal inflection points.
II. The continuity property (9) is not used in its full strength: it is enough to require that formula (9) holds for every $\boldsymbol{z}_{1} \in \mathbb{B}_{\mathcal{H}^{0}}(r)$ and $\boldsymbol{z}_{2} \in \mathbb{B}_{\mathcal{H}^{1}}\left(r_{1}\right)$. Although in most cases such a distinction is inessential, there are some situations where (9) as it is written is hard to prove.
III. In fact, we obtain the boundedness of the exponential attractor $\mathfrak{E}$ not only in $\mathcal{H}^{1}$, but also in the compactly embedded space $\mathcal{W} \Subset \mathcal{H}^{1}$. In particular, $\Pi_{2} \mathfrak{E}$ belongs to $\operatorname{dom}(T)$.
IV. A final comment on assumption (v). In certain cases, it might be difficult to prove the estimate of the compact part $K$ in $\mathcal{H}^{1}$. Actually, it is possible to relax the condition by introducing an intermediate space endowed with a norm weaker than $\mathcal{H}^{1}$, but still stronger than $\mathcal{H}^{0}$. To this end, we need to assume the existence of two reflexive Banach spaces $\hat{Y}^{1}$ and $\hat{Y}^{-1}$ such that

$$
Y^{1} \subset \hat{Y}^{1} \Subset Y^{0} \subset \hat{Y}^{-1}
$$

for which the relation

$$
\|y\|_{Y^{0}} \leq c_{0}\|y\|_{\hat{Y}^{-1}}^{\varpi}\|y\|_{\hat{Y}^{1}}^{1-\varpi}
$$

holds for every $y \in Y^{0}$ and some $c_{0}>0$ and $\varpi \in(0,1]$. Then, defining

$$
\hat{\mathcal{M}}^{1}=L_{\mu}^{2}\left(\mathbb{R}^{+} ; \hat{Y}^{1}\right)
$$

the conclusions of Theorem 5.1 remain true if $K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)$ is estimated in $X^{1} \times \hat{\mathcal{M}}^{1}$ and $\boldsymbol{w}(t)$ in $\hat{Y}^{-1}$, respectively.
6. Exponential attractors: The minimal state framework. The abstract Theorem 5.1 can be given an equivalent formulation for the semigroup $\hat{S}(t)$ acting on the extended state space $\mathcal{V}^{0}$. However, once the existence of an exponential attractor $\mathfrak{E}$ for $S(t)$ is attained, the corresponding result for $\hat{S}(t)$ can be immediately deduced under the following rather mild assumption.
Condition 6.1. For any initial datum $\hat{\boldsymbol{z}} \in \mathbb{B}_{\mathcal{V}^{0}}(r)$, denote by $\boldsymbol{x}(t)=\Pi_{1} \hat{S}(t) \hat{\boldsymbol{z}}$ the first component of the solution $\hat{S}(t) \hat{\boldsymbol{z}}$, and define

$$
\boldsymbol{\psi}^{t}(s)=\int_{0}^{\min \{t, s\}} \boldsymbol{A} \boldsymbol{x}(t-y) \mathrm{d} y
$$

Then the function $\boldsymbol{Z}(t)=\left(\boldsymbol{x}(t), \boldsymbol{\psi}^{t}\right)$ belongs to $\mathcal{H}^{0}$ for every $t \geq 0$ and

$$
\sup _{t \geq 0}\|\boldsymbol{Z}(t)\|_{\mathcal{H}^{0}} \leq Q(r)
$$

for some $Q \in \mathfrak{I}$.
Remark 5. Making use of (13), we deduce the equality

$$
\begin{equation*}
\left(\Lambda \boldsymbol{\psi}^{t}\right)(\tau)=\int_{0}^{t} \mu(\tau+s) \boldsymbol{A} \boldsymbol{x}(t-s) \mathrm{d} s \tag{15}
\end{equation*}
$$

Theorem 6.2. In addition to the general assumptions, let Condition 6.1 hold. If the semigroup $S(t): \mathcal{H}^{0} \rightarrow \mathcal{H}^{0}$ possesses an exponential attractor $\mathfrak{E}$, then the set

$$
\hat{\mathfrak{E}}=\boldsymbol{\Lambda} \mathfrak{E}
$$

is an exponential attractor for the semigroup $\hat{S}(t): \mathcal{V}^{0} \rightarrow \mathcal{V}^{0}$.

Remark 6. If $\mathfrak{E}$ is also bounded in $\mathcal{H}^{1}$, we readily infer from Lemma 4.1 the boundedness of $\hat{\mathfrak{E}}$ in $\mathcal{V}^{1}$.

Proof. By applying Lemma 4.2, and exploiting the positive invariance of $\mathfrak{E}$, we get at once

$$
\hat{S}(t) \hat{\mathfrak{E}}=\hat{S}(t) \boldsymbol{\Lambda} \mathfrak{E}=\boldsymbol{\Lambda} S(t) \mathfrak{E} \subset \boldsymbol{\Lambda} \mathfrak{E}=\hat{\mathfrak{E}} .
$$

Moreover, since $\boldsymbol{\Lambda}$ is Lipschitz continuous, we infer that $\hat{\mathfrak{E}}$ is compact and

$$
\operatorname{dim}_{\mathcal{V}^{0}}(\hat{\mathfrak{E}})=\operatorname{dim}_{\mathcal{V}^{0}}(\boldsymbol{\Lambda} \mathfrak{E}) \leq \operatorname{dim}_{\mathcal{H}^{0}}(\mathfrak{E})<\infty
$$

We are left to prove the existence of $\hat{\omega}>0$ and $Q \in \mathfrak{I}$ for which

$$
\begin{equation*}
\operatorname{dist}_{\mathcal{V}^{0}}\left(\hat{S}(t) \mathbb{B}_{\mathcal{V}^{0}}(r), \hat{\mathfrak{E}}\right) \leq Q(r) \mathrm{e}^{-\hat{\omega} t} \tag{16}
\end{equation*}
$$

To this end, let $r \geq 0$ be fixed, and let $\hat{\boldsymbol{z}}=\left(\boldsymbol{x}_{0}, \boldsymbol{\xi}_{0}\right) \in \mathbb{B}_{\mathcal{V}^{0}}(r)$. Along this proof, $C>0$ will denote a generic constant depending (increasingly) only on $r$. Setting then

$$
\hat{S}(t) \hat{\boldsymbol{z}}=\left(\boldsymbol{x}(t), \boldsymbol{\xi}^{t}\right)
$$

we construct the function $\boldsymbol{Z}(t)=\left(\boldsymbol{x}(t), \boldsymbol{\psi}^{t}\right)$ of Condition 6.1, which is bounded in $\mathcal{H}^{0}$ uniformly with respect to $t$, the bound depending on $r$. Finally, we introduce the function

$$
\hat{\boldsymbol{Z}}(t)=\boldsymbol{\Lambda} \boldsymbol{Z}(t)=\left(\boldsymbol{x}(t), \Lambda \boldsymbol{\psi}^{t}\right)
$$

The first step is showing that, for every $a, b \geq 0$,

$$
\begin{equation*}
\|\hat{S}(a+b) \hat{\boldsymbol{z}}-\hat{S}(a) \hat{\boldsymbol{Z}}(b)\|_{\mathcal{V}^{0}} \leq C \mathrm{e}^{\ell a-\frac{\delta \kappa}{2} b} \tag{17}
\end{equation*}
$$

for some $\ell=\ell(r) \geq 0$ and $\kappa=\kappa(r) \in(0,1]$. Indeed, from the continuous dependence estimate (12),

$$
\|\hat{S}(a+b) \hat{\boldsymbol{z}}-\hat{S}(a) \hat{\boldsymbol{Z}}(b)\|_{\mathcal{V}^{0}} \leq C \mathrm{e}^{\ell a}\|\hat{S}(b) \hat{\boldsymbol{z}}-\hat{\boldsymbol{Z}}(b)\|_{\mathcal{V}^{0}}^{\kappa}=C \mathrm{e}^{\ell a}\left\|\boldsymbol{\xi}^{b}-\Lambda \boldsymbol{\psi}^{b}\right\|_{\mathcal{S}^{0}}^{\kappa}
$$

On the other hand, exploiting the representation formula (11) and (15), we obtain

$$
\boldsymbol{\xi}^{b}(\tau)-\left(\Lambda \boldsymbol{\psi}^{b}\right)(\tau)=\boldsymbol{\xi}_{0}(b+\tau)
$$

Hence, by virtue of (6),

$$
\left\|\boldsymbol{\xi}^{b}-\Lambda \boldsymbol{\psi}^{b}\right\|_{\mathcal{S}^{0}}^{2}=\int_{b}^{\infty} \nu(\tau-b)\left\|\boldsymbol{\xi}_{0}(\tau)\right\|_{Y^{0}}^{2} \mathrm{~d} \tau \leq \Theta\left\|\boldsymbol{\xi}_{0}\right\|_{\mathcal{S}^{0}}^{2} \mathrm{e}^{-\delta b} \leq C \mathrm{e}^{-\delta b}
$$

This proves (17). Next, we show that, for every $a, b \geq 0$,

$$
\begin{equation*}
\operatorname{dist}_{\mathcal{V}^{0}}(\hat{S}(a) \hat{\boldsymbol{Z}}(b), \hat{\mathfrak{E}}) \leq C \mathrm{e}^{-\omega a} \tag{18}
\end{equation*}
$$

for some $\omega>0$. Indeed,

$$
\operatorname{dist}_{\mathcal{V}^{0}}(\hat{S}(a) \hat{\boldsymbol{Z}}(b), \hat{\mathfrak{E}})=\operatorname{dist}_{\mathcal{V}^{0}}(\boldsymbol{\Lambda} S(a) \boldsymbol{Z}(b), \boldsymbol{\Lambda} \mathfrak{E}) \leq \operatorname{dist}_{\mathcal{H}^{0}}(S(a) \boldsymbol{Z}(b), \mathfrak{E})
$$

and since $\boldsymbol{Z}(b)$ is uniformly bounded in $\mathcal{H}^{0}$, the desired conclusion follows from the fact that $\mathfrak{E}$ is an exponential attractor for $S(t)$. At this point, writing $t=a+b$ and making use of (17)-(18), we are led to

$$
\begin{aligned}
\operatorname{dist}_{\mathcal{V}^{0}}(\hat{S}(t) \hat{\boldsymbol{z}}, \hat{\mathfrak{E}}) & \leq\|\hat{S}(a+b) \hat{\boldsymbol{z}}-\hat{S}(a) \hat{\boldsymbol{Z}}(b)\|_{\mathcal{V}^{0}}+\operatorname{dist}_{\mathcal{V}^{0}}(\hat{S}(a) \hat{\boldsymbol{Z}}(b), \hat{\mathfrak{E}}) \\
& \leq C\left[\mathrm{e}^{\ell a-\frac{\delta \kappa}{2} b}+\mathrm{e}^{-\omega a}\right]
\end{aligned}
$$

Choosing

$$
a=\varkappa t \quad \text { with } \quad \varkappa=\frac{\delta \kappa}{\delta \kappa+2 \omega+2 \ell},
$$

the desired conclusion (16) is reached by setting $\hat{\omega}=\omega \varkappa$.

Actually, in the proof above, the exponential rate $\hat{\omega}$ depends on $r$. However, such a dependence can be removed by means of a simple argument.

Lemma 6.3. Assume that (16) holds for $\hat{\omega}=\hat{\omega}(r)$. Then (16) holds for some $\hat{\omega}$ independent of $r$ as well.

Proof. Since the set $\hat{\mathfrak{E}}$ is bounded and attracting, it is clear that for some $r_{0}>0$ large enough the set $\mathcal{B}_{0}=\mathbb{B}_{\mathcal{V}^{0}}\left(r_{0}\right)$ is absorbing for $\hat{S}(t)$. In particular, for any fixed $r \geq 0$, there exists $t_{r} \geq 0$ such that

$$
\hat{S}(t) \mathbb{B}_{\mathcal{V}^{0}}(r) \subset \mathcal{B}_{0}, \quad \forall t \geq t_{r}
$$

Hence, setting $\hat{\omega}=\hat{\omega}\left(r_{0}\right)$, we get

$$
\operatorname{dist}_{\mathcal{V}^{0}}\left(\hat{S}(t) \mathbb{B}_{\mathcal{V}^{0}}(r), \hat{\mathfrak{E}}\right) \leq \operatorname{dist}_{\mathcal{V}^{0}}\left(\hat{S}\left(t-t_{r}\right) \mathcal{B}_{0}, \hat{\mathfrak{E}}\right) \leq C \mathrm{e}^{-\hat{\omega}\left(t-t_{r}\right)}, \quad \forall t \geq t_{r}
$$

for some $C=C\left(r_{0}\right)>0$. On the other hand, it is readily inferred from (12) that

$$
\operatorname{dist}_{\mathcal{V}^{0}}\left(\hat{S}(t) \mathbb{B}_{\mathcal{V}^{0}}(r), \hat{\mathfrak{E}}\right) \leq Q(r), \quad \forall t<t_{r},
$$

for some $Q \in \mathfrak{I}$. Collecting the two inequalities above the claim follows.
7. Application to the equation of viscoelasticity. We conclude our discussion with an application of the abstract Theorems 5.1 and 6.2 to the damped wave equation with memory arising in the theory of isothermal viscoelasticity [28].
7.1. The model. Let $\Omega \subset \mathbb{R}^{3}$ be a bounded domain with smooth boundary $\partial \Omega$. For $t>0$, we consider the equation

$$
\begin{equation*}
\partial_{t t} u-\Delta u-\int_{0}^{\infty} k(s) \Delta \partial_{t} u(t-s) \mathrm{d} s+f(u)=g \tag{19}
\end{equation*}
$$

in the unknown $u=u(x, t): \Omega \times \mathbb{R} \rightarrow \mathbb{R}$, subject to the Dirichlet boundary condition

$$
u(x, t)_{\mid x \in \partial \Omega}=0 .
$$

The variable $u$ is assumed to be known for negative times $t \leq 0$, where it need not solve the equation. Defining the kernel

$$
h(s)=k(s)+1
$$

and performing an integration by parts, (19) takes the more familiar form

$$
\partial_{t t} u-h(0) \Delta u-\int_{0}^{\infty} h^{\prime}(s) \Delta u(t-s) \mathrm{d} s+f(u)=g .
$$

Besides the general assumptions of Sec. 1, we further suppose that the memory kernel $\mu$ has no jumps ${ }^{2}$ (i.e. it is absolutely continuous on $\mathbb{R}^{+}$) and

$$
\begin{equation*}
\mu^{\prime}(s)<0, \quad \text { for a.e. } s>0 . \tag{20}
\end{equation*}
$$

Concerning the other terms, $g \in L^{2}(\Omega)$ is a time-independent external force, whereas the nonlinearity $f \in \mathcal{C}^{2}(\mathbb{R})$, with $f(0)=0$, satisfies the standard growth and dissipation conditions

$$
\begin{align*}
\left|f^{\prime \prime}(u)\right| & \leq c(1+|u|),  \tag{21}\\
\liminf _{|u| \rightarrow \infty} \frac{f(u)}{u} & >-\lambda_{1}, \tag{22}
\end{align*}
$$

[^2]where $\lambda_{1}>0$ is the first eigenvalue of the Dirichlet operator $-\Delta$. Equation (19) is a particular case of the more general family of damped wave equations with memory $[6,7,12]$
$$
\partial_{t t} u+\alpha \partial_{t} u-\beta \Delta \partial_{t} u-\Delta u-\int_{0}^{\infty} k(s) \Delta \partial_{t} u(t-s) \mathrm{d} s+f(u)=g
$$
where $\alpha, \beta \geq 0$. If either $\alpha>0$ or $\beta>0$, instantaneous damping terms are present. Instead, the case $\alpha=\beta=0$ under consideration is much more challenging, since the whole dissipation mechanism is contained in the convolution integral only, and it is substantially weaker. It is also worth mentioning that (19) can be viewed as a memory relaxation of the Kelvin-Voigt model
$$
\partial_{t t} u-\Delta u-\Delta \partial_{t} u+f(u)=g
$$
the latter being recovered in the limiting situation when $k$ collapses into the Dirac mass at $0^{+}$.

Remark 7. It is readily seen that (19) can be written in the form (1) by setting

$$
\boldsymbol{x}=\binom{u}{\partial_{t} u}, \quad \boldsymbol{A} \boldsymbol{x}=\binom{0}{-\Delta \partial_{t} u}, \quad \boldsymbol{B} \boldsymbol{x}=\binom{-\partial_{t} u}{-\Delta u+f(u)-g} .
$$

7.2. Notation. Let $A=-\Delta$ be the linear strictly positive operator on $L^{2}(\Omega)$ with domain

$$
\operatorname{dom}(A)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)
$$

For $\sigma \in \mathbb{R}$, we define the compactly nested Hilbert spaces

$$
\mathrm{H}^{\sigma}=\operatorname{dom}\left(A^{\sigma / 2}\right),
$$

endowed with the inner products and norms

$$
\langle\cdot, \cdot\rangle_{\sigma}=\left\langle A^{\sigma / 2} \cdot, A^{\sigma / 2} \cdot\right\rangle_{L^{2}(\Omega)}, \quad\|\cdot\|_{\sigma}=\left\|A^{\sigma / 2} \cdot\right\|_{L^{2}(\Omega)}
$$

The index $\sigma$ will be always omitted whenever zero. In particular,

$$
\mathrm{H}^{-1}=H^{-1}(\Omega), \quad \mathrm{H}=L^{2}(\Omega), \quad \mathrm{H}^{1}=H_{0}^{1}(\Omega), \quad \mathrm{H}^{2}=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)
$$

Till the end of the paper, $Q \in \mathfrak{I}$ will denote a generic function.
7.3. Translating the equation. According to Sec. 4 and Remark 7, upon defining properly the functional spaces, equation (19) translates into an ODE in the history space, as well as in the state space frameworks. To this end, with reference to Sec. 3 , for $\imath=0,1$, we set

$$
X^{\imath}=\mathrm{H}^{\imath+1} \times \mathrm{H}^{\imath}, \quad Y^{\imath}=\{0\} \times \mathrm{H}^{\imath-1}
$$

and we define the spaces $\mathcal{M}^{\imath}, \mathcal{S}^{\imath}$ and $\mathcal{H}^{\imath}, \mathcal{V}^{2}$ accordingly. Since the first component of $Y^{\imath}$ is degenerate, abusing the notation we agree to identify $Y^{\imath}$ with its second component $\mathrm{H}^{\imath-1}$, as well as the spaces $\mathcal{M}^{\imath}, \mathcal{S}^{\imath}$ with their second components, to wit,

$$
\mathcal{M}^{\imath}=L_{\mu}^{2}\left(\mathbb{R}^{+} ; \mathrm{H}^{2-1}\right), \quad \mathcal{S}^{\imath}=L_{\nu}^{2}\left(\mathbb{R}^{+} ; \mathrm{H}^{2-1}\right)
$$

Similarly, setting

$$
\boldsymbol{\eta}=\binom{0}{\eta} \quad \text { and } \quad \boldsymbol{\xi}=\binom{0}{\xi}
$$

we will write $T \eta$ and $P \xi$ in place of $T \boldsymbol{\eta}$ and $P \boldsymbol{\xi}$, respectively. Then, the concrete realizations of the abstract equations (7) and (10) read

$$
\left\{\begin{array}{l}
\partial_{t t} u+A u+\int_{0}^{\infty} \mu(s) \eta(s) \mathrm{d} s+f(u)=g  \tag{23}\\
\partial_{t} \eta=T \eta+A \partial_{t} u
\end{array}\right.
$$

and

$$
\left\{\begin{array}{l}
\partial_{t t} u+A u+\int_{0}^{\infty} \xi(\tau) \mathrm{d} \tau+f(u)=g  \tag{24}\\
\partial_{t} \xi=P \xi+\mu A \partial_{t} u
\end{array}\right.
$$

Both (23) and (24) generate strongly (in fact, jointly) continuous semigroups

$$
S(t): \mathcal{H}^{0} \rightarrow \mathcal{H}^{0} \quad \text { and } \quad \hat{S}(t): \mathcal{V}^{0} \rightarrow \mathcal{V}^{0}
$$

Moreover, the continuous dependence estimates (9) and (12) hold for $\kappa=1$ and, due to (21)-(22), for every $r \geq 0$ we have

$$
\begin{equation*}
S(t) \mathbb{B}_{\mathcal{H}^{0}}(r) \subset \mathbb{B}_{\mathcal{H}^{0}}(Q(r)) \quad \text { and } \quad \hat{S}(t) \mathbb{B}_{\mathcal{V}^{0}}(r) \subset \mathbb{B}_{\mathcal{V}^{0}}(Q(r)) \tag{25}
\end{equation*}
$$

uniformly as $t \geq 0$ (see [5, 6, 20]). In particular, given initial data $\boldsymbol{z}=\left(u_{0}, v_{0}, \eta_{0}\right) \in$ $\mathcal{H}^{0}$ and $\hat{\boldsymbol{z}}=\left(u_{0}, v_{0}, \xi_{0}\right) \in \mathcal{V}^{0}$, we have the corresponding solutions

$$
S(t) \boldsymbol{z}=\left(u(t), \partial_{t} u(t), \eta^{t}\right) \quad \text { and } \quad \hat{S}(t) \hat{\boldsymbol{z}}=\left(u(t), \partial_{t} u(t), \xi^{t}\right)
$$

and the representation formulae (8) and (11) become

$$
\eta^{t}(s)= \begin{cases}A u(t)-A u(t-s) & s \leq t \\ \eta_{0}(s-t)+A u(t)-A u_{0} & s>t\end{cases}
$$

and

$$
\xi^{t}(\tau)=\xi_{0}(t+\tau)+\int_{0}^{t} \mu(\tau+s) A \partial_{t} u(t-s) \mathrm{d} s
$$

7.4. Exponential attractors. We are now in the position to state the main result on the existence of regular exponential attractors for the equation of viscoelasticity. In fact, although (19) is perhaps the most important (and certainly the most studied) example of equation with memory, the existence of an exponential attractor has never been proved before, not even for the much simpler model where $\mu$ satisfies the less general assumption (14) and/or an additional term of the form $\alpha \partial_{t} u$, accounting for dynamical friction, is present. So far in the literature, only the issues of global attractors and convergence of single trajectories have been addressed. In particular, the existence of the global attractor and its regularity within the hypotheses of the present work has been established in [5, 20]. This somehow justifies the need of an easy-to-handle theoretical tool allowing to treat this kind of equations.

Theorem 7.1. The semigroup $S(t)$ on $\mathcal{H}^{0}$ generated by (23) possesses an exponential attractor $\mathfrak{E}$ bounded in $\mathcal{H}^{1}$.

The proof of Theorem 7.1 will be carried out in detail in the final Sec. 8 .
Remark 8. It is worth mentioning that assumptions (2) and (20) are still much weaker than the commonly used condition (14). As a matter of fact, using the techniques devised in [25], it is even possible to weaken (20). Loosely speaking, the
conclusions of the theorem remain true if (2) holds and $\mu$ is not "too flat"; to make it precise, if the flatness rate $\mathbb{F}$ of $\mu$ does not exceed $1 / 2$, where

$$
\mathbb{F}=\frac{1}{k(0)} \int_{\left\{s: \mu^{\prime}(s)=0\right\}} \mu(y) \mathrm{d} y
$$

In our case, we assumed $\mathbb{F}=0$.
As a direct consequence of Theorem 7.1, we have
Theorem 7.2. The semigroup $\hat{S}(t)$ on $\mathcal{V}^{0}$ generated by (24) possesses an exponential attractor $\hat{\mathfrak{E}}$ bounded in $\mathcal{V}^{1}$.

Proof. Owing to Theorem 7.1, we have to show that Condition 6.1 is satisfied. In which case, the desired conclusion follows from the abstract Theorem 6.2. To this aim, let

$$
\Pi_{1} \hat{S}(t) \hat{\boldsymbol{z}}=\left(u(t), \partial_{t} u(t)\right)
$$

for an arbitrary initial datum $\hat{\boldsymbol{z}}=\left(u_{0}, v_{0}, \xi_{0}\right) \in \mathbb{B}_{\mathcal{V}^{0}}(r)$. Then, the second component (the first being trivially zero) $\psi^{t}$ of $\boldsymbol{\psi}^{t}$ fulfills the equality

$$
\psi^{t}(s)= \begin{cases}A u(t)-A u(t-s) & s \leq t \\ A u(t)-A u_{0} & s>t\end{cases}
$$

Exploiting (25), we readily obtain that

$$
\|u(t)\|_{1}+\left\|\partial_{t} u(t)\right\| \leq Q(r)
$$

and

$$
\left\|\psi^{t}(s)\right\|_{-1}=\left\|u(t)-u\left((t-s)_{+}\right)\right\|_{1} \leq Q(r)
$$

which clearly imply Condition 6.1.
8. Proof of Theorem 7.1. All we need is checking the validity of conditions (i)-(v) of Theorem 5.1. To this end we set for $\sigma \in[0,1]$

$$
\mathcal{M}^{\sigma}=L_{\mu}^{2}\left(\mathbb{R}^{+} ; \mathrm{H}^{\sigma-1}\right)
$$

and we introduce the product spaces

$$
\mathcal{H}^{\sigma}=\mathrm{H}^{\sigma+1} \times \mathrm{H}^{\sigma} \times \mathcal{M}^{\sigma}
$$

In the course of the investigation, we will borrow the following known results from [20].

- The semigroup $S(t)$ possesses a bounded absorbing set ${ }^{3} \mathbb{B}_{0} \subset \mathcal{H}^{0}$.
- For every $\varepsilon>0$ there exists $C_{\varepsilon}>0$ such that

$$
\begin{equation*}
\int_{\tau}^{t}\left\|\partial_{t} u(y)\right\| \mathrm{d} y \leq \varepsilon(t-\tau)+C_{\varepsilon} \tag{26}
\end{equation*}
$$

for every initial datum $\boldsymbol{z} \in \mathbb{B}_{0}$. Here $\partial_{t} u$ denotes the second component of the solution.

[^3]Recall that $\mathbb{B}_{0} \subset \mathcal{H}^{0}$ is an absorbing set if for every $r \geq 0$ there exists an entering time $t_{r} \geq 0$ such that

$$
S(t) \mathbb{B}_{\mathcal{H}^{0}}(r) \subset \mathbb{B}_{0}, \quad \forall t \geq t_{r}
$$

In what follows, for any initial datum $\boldsymbol{z}=\left(u_{0}, v_{0}, \eta_{0}\right) \in \mathcal{H}^{0}$ and $t \geq 0$, we agree to call

$$
S(t) \boldsymbol{z}=\left(u(t), \partial_{t} u(t), \eta^{t}\right)
$$

the solution to system (23) at time $t$. The calculations that follow are understood to hold within a suitable regularization scheme; in particular, the third component of the solution belongs to dom $(T)$. We will often use without explicit mention the Young, the Hölder and the Poincaré inequalities, as well as the standard Sobolev embeddings, e.g. $\mathrm{H}^{1} \subset L^{6}(\Omega)$ and $\mathrm{H}^{2} \subset L^{\infty}(\Omega)$.
8.1. Proof of point (i). We will actually prove a more general result.

Proposition 1. For every $\sigma \in[0,1]$, there exists $R_{\sigma}>0$ with the following property: given any $r \geq 0$ there is $q_{\sigma, r} \in \mathfrak{D}$ such that

$$
\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}} \leq q_{\sigma, r}(t)+R_{\sigma}
$$

for all $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{\sigma}}(r)$.
The proof of the case $\sigma=0$ follows by combining (25) with the existence of the absorbing set $\mathbb{B}_{0}$. Instead, the case $\sigma>0$ requires some work. In order to avoid the presence of unnecessary constants, we assume without loss of generality

$$
k(0)=\int_{0}^{\infty} \mu(s) \mathrm{d} s=1
$$

Let us take for the moment

$$
\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{\sigma}}(r) \cap \mathbb{B}_{0} .
$$

Then we know from (25) that

$$
\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{0}} \leq C
$$

where, throughout this subsection, $C>0$ denotes a generic constant, possibly depending on $\mathbb{B}_{0}$. In turn, by (21),

$$
\begin{equation*}
\|f(u(t))-g\| \leq C \tag{27}
\end{equation*}
$$

We define the energy functional $E_{\sigma}=E_{\sigma}(t)$ by

$$
E_{\sigma}=\|u\|_{\sigma+1}^{2}+\left\|\partial_{t} u\right\|_{\sigma}^{2}+\|\eta\|_{\mathcal{M}^{\sigma}}^{2}+2\left\langle f(u)-g, A^{\sigma} u\right\rangle .
$$

Since by (27) and the fact that $\sigma \leq 1$

$$
2\left|\left\langle f(u)-g, A^{\sigma} u\right\rangle\right| \leq 2\|f(u)-g\|\|u\|_{2 \sigma} \leq C\|u\|_{\sigma+1}
$$

we readily obtain the controls

$$
\begin{equation*}
\frac{1}{2}\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2}-C \leq E_{\sigma}(t) \leq 2\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2}+C \tag{28}
\end{equation*}
$$

Lemma 8.1. The functional $E_{\sigma}$ fulfills the differential inequality

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} E_{\sigma}-\int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s \leq C\left\|\partial_{t} u\right\| E_{\sigma}+C \tag{29}
\end{equation*}
$$

Proof. Multiplying the first equation of (23) by $\partial_{t} u$ in $\mathrm{H}^{\sigma}$ and the second one by $\eta$ in $\mathcal{M}^{\sigma}$, we get

$$
\frac{\mathrm{d}}{\mathrm{~d} t} E_{\sigma}=2\langle T \eta, \eta\rangle_{\mathcal{M}^{\sigma}}+2\left\langle f^{\prime}(u) \partial_{t} u, A^{\sigma} u\right\rangle .
$$

The first term in the right-hand side satisfies the identity (see e.g. [20])

$$
2\langle T \eta, \eta\rangle_{\mathcal{M}^{\sigma}}=\int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s \leq 0
$$

whereas by (21) and the standard Sobolev (or Agmon if $\sigma=1$ ) inequalities we obtain

$$
\begin{aligned}
2\left\langle f^{\prime}(u) \partial_{t} u, A^{\sigma} u\right\rangle & \leq C\left(1+\|u\|_{L^{6 /(1-\sigma)}}^{2}\right)\left\|\partial_{t} u\right\|\left\|A^{\sigma} u\right\|_{L^{6 /(1+2 \sigma)}} \\
& \leq C\left(1+\|u\|_{1}\|u\|_{\sigma+1}\right)\left\|\partial_{t} u\right\|\|u\|_{\sigma+1} \\
& \leq C\left\|\partial_{t} u\right\|\|u\|_{\sigma+1}^{2}+C .
\end{aligned}
$$

On account of (28), we are finished.
To reach the desired conclusion, we have to improve (29). To this end, an additional functional is needed to reconstruct the energy. First, in order to deal with the (possible) singularity of $\mu(s)$ at zero, for any $\nu>0$ small we choose $s_{\nu}>0$ such that

$$
\int_{0}^{s_{\nu}} \mu(s) \mathrm{d} s \leq \frac{\nu}{2}
$$

and we introduce the truncated kernel

$$
\mu_{\nu}(s)=\mu\left(s_{\nu}\right) \chi_{\left(0, s_{\nu}\right]}(s)+\mu(s) \chi_{\left(s_{\nu}, \infty\right)}(s)
$$

where $\chi$ denotes the characteristic function. Besides, for any $\delta>0$ we set

$$
\begin{aligned}
\mathcal{P}_{\delta}[\eta] & =\int_{P_{\delta}} \mu(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s \\
\mathcal{N}_{\delta}[\eta] & =\int_{N_{\delta}} \mu(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s
\end{aligned}
$$

where

$$
P_{\delta}=\left\{s \in \mathbb{R}^{+}: \mu^{\prime}(s)+\delta \mu(s)>0\right\}, \quad N_{\delta}=\left\{s \in \mathbb{R}^{+}: \mu^{\prime}(s)+\delta \mu(s) \leq 0\right\}
$$

Note that

$$
\begin{equation*}
\mathcal{N}_{\delta}[\eta] \leq-\frac{1}{\delta} \int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s \tag{30}
\end{equation*}
$$

and

$$
\mathcal{P}_{\delta}[\eta]+\mathcal{N}_{\delta}[\eta]=\|\eta\|_{\mathcal{M}^{\sigma}}^{2}
$$

Finally, we define the functional $\Phi=\Phi(t)$ as

$$
\begin{aligned}
\Phi= & -\int_{0}^{\infty} \mu_{\nu}(s)\left\langle\partial_{t} u, \eta(s)\right\rangle_{\sigma-1} \mathrm{~d} s+(1-2 \nu)\left\langle\partial_{t} u, u\right\rangle_{\sigma} \\
& +\int_{0}^{\infty}\left(\int_{s}^{\infty} \mu(y) \chi_{P_{\delta}}(y) \mathrm{d} y\right)\|\eta(s)-A u\|_{\sigma-1}^{2} \mathrm{~d} s .
\end{aligned}
$$

Since, as shown in [20], condition (2) is equivalent to

$$
k(s) \leq D \mu(s)
$$

for some $D>0$, it is immediate to ascertain that

$$
\begin{equation*}
|\Phi(t)| \leq C\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2} \tag{31}
\end{equation*}
$$

Collecting Lemmas 5.3 and 5.4 in [20], and taking advantage of (30), we obtain the following lemma, ${ }^{4}$ whose almost immediate proof is left to the reader.

Lemma 8.2. For any $\nu, \delta>0$ small, there exists a positive constant $K=K(\nu, \delta)$ such that the following inequality holds:

$$
\begin{aligned}
& \frac{\mathrm{d}}{\mathrm{~d} t} \Phi+\frac{1}{4}\|u\|_{\sigma+1}^{2}+\nu\left\|\partial_{t} u\right\|_{\sigma}^{2}+\frac{1}{4}\|\eta\|_{\mathcal{M}^{\sigma}}^{2} \\
& \leq-K \int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s+K\|f(u)-g\|^{2}
\end{aligned}
$$

Accordingly, up to fixing $\nu, \delta>0$ small enough and recalling (27), we conclude that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \Phi+\nu\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2} \leq-C \int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s+C \tag{32}
\end{equation*}
$$

Conclusion of the proof of Proposition 1. At this point, for $\varepsilon>0$ we introduce the further functional $\Gamma=\Gamma(t)$ as

$$
\Gamma=E_{\sigma}+\varepsilon \Phi .
$$

By (28) and (31) it is apparent that

$$
\begin{equation*}
\frac{1}{4}\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2}-C \leq \Gamma(t) \leq 4\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2}+C \tag{33}
\end{equation*}
$$

for any $\varepsilon>0$ small enough. By collecting (29) and (32) we obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \Gamma+\nu \varepsilon\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{\sigma}}^{2} \leq(1-C \varepsilon) \int_{0}^{\infty} \mu^{\prime}(s)\|\eta(s)\|_{\sigma-1}^{2} \mathrm{~d} s+C\left\|\partial_{t} u\right\| E_{\sigma}+C
$$

Choosing then $\varepsilon$ sufficiently small, and recalling (28) and (33), we end up with the differential inequality

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \Gamma+2 \gamma \Gamma \leq C\left\|\partial_{t} u\right\| \Gamma+C
$$

for some $\gamma>0$. Owing to the dissipation integral (26), we can apply a modified version of the Gronwall lemma (see e.g. Lemma 2.1 in [6]) to get

$$
\Gamma(t) \leq C|\Gamma(0)| \mathrm{e}^{-\gamma t}+C
$$

A final use of (33) completes the proof when $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{\sigma}}(r) \cap \mathbb{B}_{0}$.
As far as the general case is concerned, we observe that there exists $t_{r} \geq 0$ such that $S\left(t_{r}\right) \mathbb{B}_{\mathcal{H}^{\sigma}}(r) \subset \mathbb{B}_{0}$. It is then enough to show that

$$
\begin{equation*}
S(t) \mathbb{B}_{\mathcal{H}^{\sigma}}(r) \subset \mathbb{B}_{\mathcal{H}^{\sigma}}(Q(r)), \quad \forall t \leq t_{r} \tag{34}
\end{equation*}
$$

Indeed, once this is known, we can write (for $t \geq t_{r}$ )

$$
S(t) \boldsymbol{z}=S\left(t-t_{r}\right) \tilde{\boldsymbol{z}} \quad \text { with } \quad \tilde{\boldsymbol{z}}=S\left(t_{r}\right) \boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{\sigma}}(Q(r)) \cap \mathbb{B}_{0}
$$

and the claim follows by the previous step. The proof of (34) can be done by merely applying the Gronwall lemma to (29) on the time interval $\left[0, t_{r}\right]$, the only difference being that now the constant $C$ will depend on $r$.

[^4]8.2. Proof of point (ii). Let $r \geq 0$ be given. According to [20] (see Lemmas 4.3 and 4.4 therein), for every $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{0}}(r)$ the solution $S(t) \boldsymbol{z}$ splits into the sum
$$
S(t) \boldsymbol{z}=L(t) \boldsymbol{z}+K(t) \boldsymbol{z}
$$
where
$$
\|L(t) \boldsymbol{z}\|_{\mathcal{H}^{0}} \leq Q(r) \mathrm{e}^{-\omega t} \quad \text { and } \quad\|K(t) \boldsymbol{z}\|_{\mathcal{H}^{1 / 3}} \leq Q(r)
$$

Here, $\omega>0$ is actually independent of $r$. In particular, the (bounded) absorbing set $\mathbb{B}_{0}$ is exponentially attracted by a ball $\mathbb{B}_{1 / 3}$ of $\mathcal{H}^{1 / 3}$. Besides, by Proposition 1 with $\sigma=1 / 3$, we know that $S(t) \mathbb{B}_{1 / 3}$ remains uniformly bounded (with respect to $t$ ) in $\mathcal{H}^{1 / 3}$. Hence, Lemma 4.5 in [20] allows us to draw the uniform bound

$$
\sup _{\boldsymbol{z} \in \mathbb{B}_{1 / 3}}\|K(t) \boldsymbol{z}\|_{\mathcal{H}^{1}} \leq C
$$

for some $C>0$ depending only on $\mathbb{B}_{1 / 3}$. This, together with the decay of $L(t) \boldsymbol{z}$, tell that $\mathbb{B}_{1 / 3}$ is exponentially attracted by a ball $\mathbb{B}_{1}$ of $\mathcal{H}^{1}$. By the transitivity property of exponential attraction [16], which applies since we have the continuity (9), we infer that $\mathbb{B}_{0}$ is exponentially attracted by $\mathbb{B}_{1}$. At this point, since $\mathbb{B}_{0}$ is absorbing and (25) holds, it is a standard matter to conclude that $\mathbb{B}_{1}$ exponentially attracts every bounded subset of $\mathcal{H}^{0}$.
8.3. Proof of point (iii). For $\boldsymbol{x}=(u, v)$, it is immediate to see that

$$
\|\boldsymbol{A} \boldsymbol{x}\|_{Y^{0}}=\|A v\|_{-1}=\|v\|_{1} \leq\|\boldsymbol{x}\|_{X^{1}},
$$

yielding the desired bound.
8.4. Proof of point (iv). If $\boldsymbol{z} \in \mathbb{B}_{\mathcal{H}^{1}}(r)$, by the previous point (i) we get

$$
\|S(t) \boldsymbol{z}\|_{\mathcal{H}^{1}} \leq Q(r) \quad \Rightarrow \quad\left\|\partial_{t} u(t)\right\|_{1} \leq Q(r)
$$

Recalling (21), we read from the first equation of (23) the further bound

$$
\left\|\partial_{t t} u(t)\right\| \leq Q(r)
$$

Therefore, (iv) holds with $p=\infty$.
8.5. Proof of point (v). Let $r \geq 0$ be fixed. Along this proof, $C>0$ and $Q \in \mathfrak{I}$ will denote a generic constant and function, respectively, both possibly depending on $r$. For $\boldsymbol{z}_{1}, \boldsymbol{z}_{2} \in \mathbb{B}_{\mathcal{H}^{1}}(r)$ arbitrarily chosen, we deduce from point (i) the bound

$$
\begin{equation*}
\left\|S(t) \boldsymbol{z}_{\imath}\right\|_{\mathcal{H}^{1}}=\left\|\left(u_{\imath}(t), \partial_{t} u_{\imath}(t), \eta_{\imath}^{t}\right)\right\|_{\mathcal{H}^{1}} \leq Q(r) \tag{35}
\end{equation*}
$$

We decompose the difference

$$
S(t) \boldsymbol{z}_{1}-S(t) \boldsymbol{z}_{2}=\left(\bar{u}(t), \partial_{t} \bar{u}(t), \bar{\eta}^{t}\right)
$$

as follows:

$$
S(t) \boldsymbol{z}_{1}-S(t) \boldsymbol{z}_{2}=L\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)+K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)
$$

where

$$
L\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=\left(v(t), \partial_{t} v(t), \xi^{t}\right) \quad \text { and } \quad K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=\left(w(t), \partial_{t} w(t), \psi^{t}\right)
$$

fulfill the systems

$$
\left\{\begin{array}{l}
\partial_{t t} v+A v+\int_{0}^{\infty} \mu(s) \xi(s) \mathrm{d} s=0 \\
\partial_{t} \xi=T \xi+A \partial_{t} v
\end{array}\right.
$$

and

$$
\left\{\begin{array}{l}
\partial_{t t} w+A w+\int_{0}^{\infty} \mu(s) \psi(s) \mathrm{d} s+f\left(u_{1}\right)-f\left(u_{2}\right)=0 \\
\partial_{t} \psi=T \psi+A \partial_{t} w
\end{array}\right.
$$

with initial data

$$
L\left(0, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=\boldsymbol{z}_{1}-\boldsymbol{z}_{2} \quad \text { and } \quad K\left(0, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=0 .
$$

As shown in $[25,26]$, the first (linear) system generates an exponentially stable contraction semigroup on $\mathcal{H}^{0}$. Thus,

$$
\begin{equation*}
\left\|L\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{0}} \leq B \mathrm{e}^{-\omega t}\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}} \tag{36}
\end{equation*}
$$

for some $B \geq 1$ and $\omega>0$. Concerning the second system, we define the energy functionals

$$
E_{K}(t)=\left\|K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{1}}^{2}=\|w(t)\|_{2}^{2}+\left\|\partial_{t} w(t)\right\|_{1}^{2}+\left\|\psi^{t}\right\|_{\mathcal{M}^{1}}^{2}
$$

and

$$
\Lambda(t)=E_{K}(t)+2\left\langle f\left(u_{1}(t)\right)-f\left(u_{2}(t)\right), w(t)\right\rangle_{1} .
$$

Owing to (21) and (35), we draw

$$
2\left|\left\langle f\left(u_{1}\right)-f\left(u_{2}\right), w\right\rangle_{1}\right| \leq C\|\bar{u}\|_{1}\|w\|_{2} \leq \frac{1}{2}\|w\|_{2}^{2}+C\|\bar{u}\|_{1}^{2} \leq \frac{1}{2} E_{K}+C\|\bar{u}\|_{1}^{2}
$$

Since by (9) (which holds for $\kappa=1$ )

$$
\|\bar{u}(t)\|_{1}^{2} \leq Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2}
$$

we conclude that

$$
\begin{equation*}
\Lambda(t) \geq \frac{1}{2} E_{K}(t)-Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2} \tag{37}
\end{equation*}
$$

Besides, by direct calculations,

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} t} \Lambda & =2\left\langle f^{\prime}\left(u_{1}\right) \partial_{t} \bar{u}, w\right\rangle_{1}+2\left\langle\left(f^{\prime}\left(u_{1}\right)-f^{\prime}\left(u_{2}\right)\right) \partial_{t} u_{2}, w\right\rangle_{1}+2\langle T \psi, \psi\rangle_{\mathcal{M}^{1}} \\
& \leq 2\left\langle f^{\prime}\left(u_{1}\right) \partial_{t} \bar{u}, w\right\rangle_{1}+2\left\langle\left(f^{\prime}\left(u_{1}\right)-f^{\prime}\left(u_{2}\right)\right) \partial_{t} u_{2}, w\right\rangle_{1}
\end{aligned}
$$

Appealing once again to (21) and (35),

$$
2\left\langle f^{\prime}\left(u_{1}\right) \partial_{t} \bar{u}, w\right\rangle_{1} \leq C\left\|\partial_{t} \bar{u}\right\|\|w\|_{2}
$$

and

$$
2\left\langle\left(f^{\prime}\left(u_{1}\right)-f^{\prime}\left(u_{2}\right)\right) \partial_{t} u_{2}, w\right\rangle_{1} \leq C\|\bar{u}\|_{L^{6}}\left\|\partial_{t} u_{2}\right\|_{L^{3}}\|w\|_{2} \leq C\|\bar{u}\|_{1}\|w\|_{2}
$$

Therefore, making use of (36) and (37),

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} t} \Lambda(t) & \leq C\|w(t)\|_{2}\left(\|\bar{u}(t)\|_{1}+\left\|\partial_{t} \bar{u}(t)\right\|\right) \\
& \leq C\|w(t)\|_{2}\left(\|w(t)\|_{1}+\left\|\partial_{t} w(t)\right\|+\|v(t)\|_{1}+\left\|\partial_{t} v(t)\right\|\right) \\
& \leq C\left(E_{K}(t)+\sqrt{E_{K}(t)}\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}\right) \\
& \leq C\left(E_{K}(t)+\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2}\right) \\
& \leq C \Lambda(t)+Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2} .
\end{aligned}
$$

Observing that $\Lambda(0)=0$, an application of the Gronwall Lemma yields

$$
\Lambda(t) \leq Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2}
$$

and a further use of (37) leads to

$$
\left\|K\left(t, \boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)\right\|_{\mathcal{H}^{1}}^{2}=E_{K}(t) \leq Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}^{2}
$$

Accordingly,

$$
\left\|A \partial_{t} w(t)\right\|_{-1}=\left\|\partial_{t} w(t)\right\|_{1} \leq \sqrt{E_{K}(t)} \leq Q(t)\left\|\boldsymbol{z}_{1}-\boldsymbol{z}_{2}\right\|_{\mathcal{H}^{0}}
$$

Collecting (36) and the last two estimates, the claim follows.

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[^1]:    ${ }^{1}$ Actually, we are using a well-known generalization of the result in [4, 14], originally stated with the decomposition $S \boldsymbol{z}=L \boldsymbol{z}+K \boldsymbol{z}$. In turn, $L\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=L \boldsymbol{z}_{1}-L \boldsymbol{z}_{2}$ and $K\left(\boldsymbol{z}_{1}, \boldsymbol{z}_{2}\right)=K \boldsymbol{z}_{1}-K \boldsymbol{z}_{2}$. Indeed, a closer look at the proofs of $[4,14]$ shows that it suffices to split the difference $S \boldsymbol{z}_{1}-S \boldsymbol{z}_{2}$ as in (v).

[^2]:    ${ }^{2}$ The request that $\mu$ have no jumps is actually made only to simplify the notation. Indeed, in the presence of jumps, no significant changes are needed in the forthcoming proofs.

[^3]:    ${ }^{3}$ Actually, the proof of the existence of an absorbing set in [20] is indirect, being a consequence of the existence of the global attractor, which is attained by means of gradient system techniques.

[^4]:    ${ }^{4}$ Actually, Lemmas 5.3 and 5.4 in [20] are proved for $\sigma=0$, but the proofs are in fact the same for any $\sigma \in[0,1]$ (cf. [20, Remark 5.5]).

