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AN EULER-BERNOULLI BEAM WITH NONLINEAR DAMPING AND A NONLINEAR SPRING AT THE TIP

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ABSTRACT. We study the asymptotic behaviour for a system consisting of a clamped flexible beam that carries a tip payload, which is attached to a nonlinear damper and a nonlinear spring at its end. Characterizing the ω -limit sets of the trajectories, we give a necessary condition under which the system is asymptotically stable. In the case when this condition is not satisfied, we show that the beam deflection approaches a non-decaying time-periodic solution.

1. Introduction. This article considers an Euler-Bernoulli beam where one end is clamped, and the free end holds a rigid tip mass. Models of this form play a fundamental role in many mechanical systems and thus occur in many applications such as flexible robot arms, helicopter rotor blades, spacecraft antennae, airplane wings, high-rise buildings, etc. An important issue is the suppression of vibrations, since undesired oscillations can reduce the performance of the system, or worse, result in damage to the structure. For this reason, the Euler-Bernoulli beam is often coupled with a boundary control, which acts on the tip and is used to dissipate the vibration. Frequently, the boundary control is realized as a suspension system, consisting typically of springs and dampers.

In the last four decades, considerable attention has been paid to the stability analysis of such systems in the literature, see e.g. [15, 12, 11, 17]. Most results deal with the situation in which the control is linear, thereby obtaining linear boundary conditions. In general, the respective stability analysis uses results from linear functional analysis. However, the generalization to nonlinear boundary conditions is not straightforward in most cases, since the linear techniques do not apply in this situation any more. Up to the knowledge of the authors the only models considered in the literature with nonlinearities at the boundary do not have a rigid body attached to the tip (see for example [4, 5, 6]). A major purpose of this article is to make a first step towards closing this gap.

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In this article we investigate an Euler-Bernoulli beam which is clamped at one end (see Figure 1). The free end holds a tip mass, whose mass m and moment of inertia J are both positive. The controller acting on the tip consists of a spring and a damper, both nonlinear. This is a convenient nonlinear model to start the analysis with, and to investigate the asymptotic behaviour. While being a toy model, the authors are confident that the analysis presented in this article can easily be extended to more complex models of this kind.

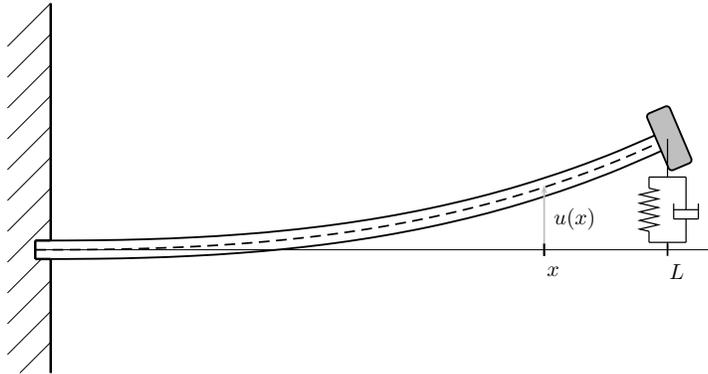


FIGURE 1. Clamped beam with tip mass, coupled to a spring and damper (both nonlinear)

We consider an Euler-Bernoulli beam satisfying a linear PDE with high order nonlinear boundary conditions. In order to make the system accessible for analysis it is a common strategy to rewrite it as a nonlinear evolution equation in an appropriate (infinite-dimensional) Hilbert space \mathcal{H} . In general, proving that every mild solution tends to zero as time goes to infinity consists of two steps, namely showing the precompactness of the trajectories and proving that the only possible limit is the zero solution. In the linear case, verifying the precompactness is straightforward by showing that the resolvent of the system operator is compact, see Remark 4.5 below. For the nonlinear case, the inspection of the precompactness property is more complex. The most commonly used criteria for the precompactness of trajectories can be found in [8, 20, 19, 25], and further generalizations in [7, 24]. They all split the system operator into the sum of two operators $A + \mathcal{N}$ (A being its linear, and \mathcal{N} its nonlinear part) and infer precompactness under the following conditions. In [8] A is required to be m -dissipative and \mathcal{N} applied to a trajectory is L^1 in time. In [19] the requirement on \mathcal{N} is loosened by just assuming uniform local integrability of \mathcal{N} applied to a trajectory. However, the linear semigroup e^{tA} needs additionally to be compact in order to still ensure precompactness. Finally, in [25] \mathcal{N} just needs to map into a compact set, but A needs to generate an exponentially stable linear C_0 -semigroup. These strategies have successfully been applied in the literature to the Euler-Bernoulli beam without tip payload and with nonlinear boundary control: in [5] the precompactness of the trajectories follows directly from the m -dissipativity of the system operator, and in [4] from the L^1 -integrability of the nonlinearity.

In contrast to the mentioned literature, the nonlinear boundary conditions considered in this article do not fall into any of these sets of assumptions. In our case

A shall be m -dissipative, but not compact and it does not generate an exponentially stable semigroup. On the other hand, \mathcal{N} apparently does not satisfy strong assumptions either, for it is compact, but we can not guarantee L^1 -integrability. Thus the properties of our system operator are too weak in order to apply the mentioned standard results. However, we are still able to prove precompactness of the trajectories in a novel way. Hence, this article enriches the available “tool box” for such evolution problems.

In this article we show that for the Euler-Bernoulli beam with tip mass, coupled to a nonlinear spring and a nonlinear damper, all trajectories that are C^1 in time are precompact. Furthermore, for almost all values of the moment of inertia $J > 0$ the trajectories tend to zero as time goes to infinity. Interestingly we find that, for countably many values of J , the trajectories tend to a time-periodic solution. For given initial conditions we are able to characterize this asymptotic limit explicitly, including its phase. Such periodic limiting orbit appears when the (linear) beam equation has an eigenfunction with a node at the free end (i.e. some $u_n(L) = 0$). Then, the controller at the tip is inactive for all time.

A possible application of the method developed here is the nonlinear extension of the linear theory in [2], describing a model for a flexible microgripper used for DNA manipulation (the DNA-bundle model consists of a damper, spring and a load). Studying the stability of the system, when nonlinear phenomena for the controller and DNA-bundle are included, is a goal for future research set by the authors.

The paper is organized as follows. In Section 2 the equations of motion are derived for the system consisting of the Euler-Bernoulli beam with tip mass, connected to a nonlinear spring and damper. Next, it is shown that the energy functional is an appropriate Lyapunov function for the system. Section 3 is concerned with the formulation of the problem in an appropriate functional analytical setting and the investigation of existence and uniqueness of the corresponding solutions. In Section 4 we prove precompactness of the trajectories for all initial conditions lying in a dense subset of the underlying Hilbert space. Section 5 deals with the characterization of possible ω -limit sets, proving that any classical solution tends either to zero or to a periodic solution, depending on the prescribed value J .

2. Preliminaries and derivation of the model. For a function $u(t, x)$, $t \geq 0$, $x \in [0, L]$ for some $L > 0$, we use the notation u_t for the derivative with respect to the time variable t , and we write u' for the x -derivative. Higher order x -derivatives are also denoted by roman superscripts. Whenever it is clear from the context, we omit the time variable in the notation and write for example $u(L) \equiv u(t, L)$ and $u(0) \equiv u(t, 0)$. If not stated otherwise all functions occurring in this article are considered to be real valued, and only real valued solutions are sought. Therefore, in addition to the Hilbert spaces $L^2(0, L)$ and $H^k(0, L)$, which are understood to consist of complex valued functions, we also need $L^2_{\mathbb{R}}(0, L) := \{f \in L^2(0, L) : f : [0, L] \rightarrow \mathbb{R}\}$, and analogously we define $H^k_{\mathbb{R}}(0, L)$.

A linear operator A is a closed, linear map $A : \mathcal{X} \rightarrow \mathcal{X}$, where \mathcal{X} is a suitable real or complex Hilbert space. The operator A is defined on the domain $D(A)$ which needs to be dense in \mathcal{X} . The range of A is $\text{ran } A \subset \mathcal{X}$. A closed linear subset X of a Hilbert space \mathcal{X} is called A -invariant if $X \cap D(A) \subset X$ is dense and $\text{ran } A|_X \subset X$.

Throughout this article C denotes a positive constant, not necessarily always the same.

For the derivation of the model we follow [10] and [14], whereby we assume that the beam satisfies the Euler-Bernoulli assumption. We assume that the beam has uniform mass per length $\rho > 0$ and length L . The beam is parametrized with $x \in [0, L]$, and is described by its deviation $u(t, x)$ from the horizontal (as depicted in Figure 1). The constant bending stiffness is $\Lambda > 0$, and the tension is assumed to be zero. At the tip of the beam there is a payload of mass $m > 0$, which has the moment of inertia $J > 0$. We neglect friction of any kind. Only two external forces are assumed to act on the beam, both on the tip, perpendicular to the resting position $u \equiv 0$. The first comes from a nonlinear spring attached to the tip, producing the restoring force $-k_1(u(t, L))$. The second force is due to a nonlinear damping, and is given by $-k_2(u_t(t, L))$. We assume $k_1 \in C^2(\mathbb{R})$ and $k_2 \in C^2(\mathbb{R})$, and

$$\int_0^z k_1(s) ds \geq 0, \quad \forall z \in \mathbb{R}, \quad (2.1)$$

$$k_2'(z) \geq 0, \quad k_2(0) = 0, \quad \forall z \in \mathbb{R}. \quad (2.2)$$

Furthermore, assume that

$$|k_2(z)| \geq Kz^2, \quad \forall z \in (-\delta, \delta), \quad (2.3)$$

for some positive constant $K > 0$ and $\delta > 0$ small. Notice that (2.1) implies $k_1(0) = 0$, and (2.2) together with (2.3) imply that $k_2(z) = 0$ iff $z = 0$.

The equations of motion can be derived according to Hamilton's principle, i.e. they are the Euler-Lagrange equations corresponding to the action functional. In our model the kinetic energy E_k and the potential (strain) energy E_p are

$$E_k = \frac{\rho}{2} \int_0^L u_t^2 dx + \frac{m}{2} u_t(L)^2 + \frac{J}{2} u_t'(L)^2, \quad E_p = \frac{\Lambda}{2} \int_0^L (u'')^2 dx.$$

Additionally, we have the virtual work δW coming from the external forces:

$$\delta W = -k_1(u(L))\delta u(L) - k_2(u_t(L))\delta u(L).$$

Taking into account the boundary conditions $u(0) = u'(0) = 0$ of the clamped end we find that, according to Hamilton's principle, u satisfies:

$$\rho u_{tt}(t, x) + \Lambda u^{IV}(t, x) = 0, \quad (2.4a)$$

$$u(t, 0) = u'(t, 0) = 0, \quad (2.4b)$$

$$-\Lambda u'''(t, L) + m u_{tt}(t, L) = -k_1(u(t, L)) - k_2(u_t(t, L)), \quad (2.4c)$$

$$\Lambda u''(t, L) + J u_t'(t, L) = 0, \quad (2.4d)$$

where $(t, x) \in (0, \infty) \times (0, L)$.

Finally we derive a candidate for a Lyapunov function, guided by the expectation that the total energy of the beam will decrease in time, mainly due to the damping. The total energy of the system is given by $E_{\text{tot}} = E_k + E_p + E_s$, where $E_s := \int_0^{u(L)} k_1(s) ds$ represents the potential energy stored in the nonlinear spring. Now (2.1) ensures that this integral always stays non-negative. We compute the time

derivative of the total energy, using the Euler-Lagrange equations (2.4):

$$\begin{aligned}
\frac{d}{dt}E_{\text{tot}} &= \Lambda \int_0^L u'' u_t'' dx + \rho \int_0^L u_{tt} u_t dx + m u_{tt}(L) u_t(L) + J u_{tt}'(L) u_t'(L) \\
&\quad + k_1(u(L)) u_t(L) \\
&= \Lambda \int_0^L u^{\text{IV}} u_t dx + \Lambda u'' u_t'|_0^L - \Lambda u''' u_t|_0^L + \rho \int_0^L u_{tt} u_t dx \\
&\quad + m u_{tt}(L) u_t(L) + J u_{tt}'(L) u_t'(L) + k_1(u(L)) u_t(L) \\
&= \Lambda u''(L) u_t'(L) - \Lambda u'''(L) u_t(L) + m u_{tt}(L) u_t(L) + J u_{tt}'(L) u_t'(L) \\
&\quad + k_1(u(L)) u_t(L) \\
&= -k_2(u_t(L)) u_t(L).
\end{aligned}$$

Due to (2.2) this last line is always non-positive. So E_{tot} is a candidate for a Lyapunov function:

$$V(u) := \frac{\Lambda}{2} \int_0^L (u'')^2 dx + \frac{\rho}{2} \int_0^L u_t^2 dx + \frac{m}{2} u_t(L)^2 + \frac{J}{2} u_t'(L)^2 + \int_0^{u(L)} k_1(s) ds, \quad (2.5)$$

and it is non-negative. According to the previous calculation its derivative along classical solutions of (2.4) satisfies:

$$\frac{d}{dt}V(u(t)) = -k_2(u_t(t, L)) u_t(t, L) \leq 0. \quad (2.6)$$

3. Formulation as an evolution equation. The aim of this section is to show that for sufficiently regular initial conditions $u(0, x) = u_0(x)$ and $u_t(0, x) = v_0(x)$ the system (2.4) has a unique (mild) solution $u(t, x)$. First, we introduce the standard setting for the Euler-Bernoulli beam with a tip payload (see [14], [17]). To this end we define the following real Hilbert space:

$$\mathcal{H} := \{y = [u, v, \xi, \psi]^\top : u \in \tilde{H}_{0, \mathbb{R}}^2(0, L), v \in L_{\mathbb{R}}^2(0, L), \xi, \psi \in \mathbb{R}\},$$

where $\tilde{H}_{0, \mathbb{R}}^n(0, L) := \{f \in H_{\mathbb{R}}^n(0, L) : f(0) = f'(0) = 0\}$ for $n \geq 2$. The space \mathcal{H} is equipped with the inner product

$$\langle y_1, y_2 \rangle_{\mathcal{H}} := \frac{\Lambda}{2} \int_0^L u_1'' u_2'' dx + \frac{\rho}{2} \int_0^L v_1 v_2 dx + \frac{1}{2J} \xi_1 \xi_2 + \frac{1}{2m} \psi_1 \psi_2, \quad \forall y_1, y_2 \in \mathcal{H}.$$

Next we consider the following linear operator on \mathcal{H} :

$$A : y \mapsto \begin{bmatrix} v \\ -\frac{\Lambda}{\rho} u^{\text{IV}} \\ -\Lambda u''(L) \\ \Lambda u'''(L) \end{bmatrix} \quad (3.1)$$

on the dense domain

$$D(A) := \{y \in \mathcal{H} : u \in \tilde{H}_{0, \mathbb{R}}^4(0, L), v \in \tilde{H}_{0, \mathbb{R}}^2(0, L), \xi = Jv'(L), \psi = mv(L)\}.$$

Furthermore, we define the bounded nonlinear operator \mathcal{N} on \mathcal{H} :

$$\mathcal{N} : y \mapsto \begin{bmatrix} 0 \\ 0 \\ 0 \\ -k_1(u(L)) - k_2\left(\frac{\psi}{m}\right) \end{bmatrix}.$$

Finally we introduce the nonlinear operator $\mathcal{A} := A + \mathcal{N}$ on the domain $D(\mathcal{A}) = D(A)$. With this notation the system (2.4) can be written formally as the following nonlinear evolution equation in \mathcal{H} :

$$y_t = \mathcal{A}y, \quad (3.2a)$$

$$y(0) = y_0, \quad (3.2b)$$

for some initial condition $y_0 \in \mathcal{H}$.

A function $y(t)$ is a classical solution of (3.2) on $(0, T)$ if $y \in C^1((0, T); \mathcal{H}) \cap C([0, T]; \mathcal{H})$ and for all $t \in (0, T)$ there holds $y(t) \in D(\mathcal{A})$ and (3.2). A continuous function $y \in C([0, T]; \mathcal{H})$ that satisfies the Duhamel formula

$$y(t) = e^{tA}y_0 + \int_0^t e^{(t-s)A}\mathcal{N}y(s) ds, \quad t \in (0, T), \quad (3.3)$$

is called a mild solution of (3.2), see [21].

Lemma 3.1. *The linear operator A generates a C_0 -semigroup $(e^{tA})_{t \geq 0}$ of unitary operators in \mathcal{H} .*

For the proof see the Lemma A.2 in the appendix. For the following analysis we need to properly define a Lyapunov function on \mathcal{H} . In the previous section we obtained a candidate by (2.5), which was defined along classical solutions $u(t)$ of (2.4). For $y \in D(\mathcal{A})$ this can equivalently be written as

$$V(y) = \frac{\Lambda}{2} \int_0^L (u'')^2 dx + \frac{\rho}{2} \int_0^L v^2 dx + \frac{1}{2m} \psi^2 + \frac{1}{2J} \xi^2 + \int_0^{u(L)} k_1(s) ds.$$

This is then defined for all $y \in \mathcal{H}$, and according to (2.6) there holds

$$\frac{d}{dt} V(y(t)) = -k_2 \left(\frac{\psi(t)}{m} \right) \frac{\psi(t)}{m} \leq 0, \quad (3.4)$$

along all classical solutions $y(t)$. The following two results are easily verified:

Lemma 3.2. *The function $V : \mathcal{H} \rightarrow \mathbb{R}$ is continuous with respect to the norm $\|\cdot\|_{\mathcal{H}}$.*

Lemma 3.3. *For any $Y \subset \mathcal{H}$ we have:*

$$\sup \{V(y) : y \in Y\} < \infty \quad \Leftrightarrow \quad \sup \{\|y\|_{\mathcal{H}} : y \in Y\} < \infty$$

Proposition 3.4. *For every $y_0 \in \mathcal{H}$ there exists a unique mild solution $y : [0, T_{\max}(y_0)) \rightarrow \mathcal{H}$, where $T_{\max}(y_0)$ is the maximal time interval for which the solution exists. If $T_{\max}(y_0) < \infty$ then a blow-up occurs, i.e.*

$$\lim_{t \nearrow T_{\max}} \|y(t)\|_{\mathcal{H}} = \infty.$$

Proof. Due to the assumptions made on k_1, k_2 it follows that \mathcal{N} is continuously differentiable on \mathcal{H} , and thus locally Lipschitz continuous. Furthermore, A generates a C_0 -semigroup. Hence, according to Theorem 6.1.4 in [21] a unique mild solution exists on $[0, T_{\max}(y_0))$, for some maximal $0 < T_{\max}(y_0) \leq \infty$. Moreover, if $T_{\max} < \infty$ then $\lim_{t \nearrow T_{\max}} \|y(t)\|_{\mathcal{H}} = \infty$. \square

Lemma 3.5. *If $y_0 \in D(\mathcal{A})$ then the corresponding mild solution $y(t)$ is a classical solution. Furthermore $y(t)$ is a global solution, i.e. $T_{\max}(y_0) = \infty$.*

Proof. Since \mathcal{N} is continuously differentiable, Theorem 6.1.5 in [21] implies that $y(t)$ is a classical solution. Therefore (2.6) holds and implies:

$$V(y(t)) \leq V(y_0), \quad \forall t \in [0, T_{\max}).$$

Thus, according to Lemma 3.3 the norm $\|y(t)\|_{\mathcal{H}}$ stays uniformly bounded on $[0, T_{\max})$. Consequently, no blowup occurs and $T_{\max} = \infty$. \square

The following result is a consequence of Proposition 4.3.7 of [3]:

Proposition 3.6. *Let $y : [0, T) \rightarrow \mathcal{H}$ be a mild solution of (3.2) for some $y_0 \in \mathcal{H}$, and $0 < T \leq \infty$. Also, let $\{y_{n,0}\}_{n \in \mathbb{N}} \subset D(\mathcal{A})$ be such that $y_{n,0} \rightarrow y_0$ in \mathcal{H} . For every $n \in \mathbb{N}$, denote by $y_n(t)$ the (global) classical solution of (3.2) to the initial value $y_{n,0}$. Then $y_n \rightarrow y$ in $C([0, T]; \mathcal{H})$ as $n \rightarrow \infty$.*

Theorem 3.7. *For every $y_0 \in \mathcal{H}$ the initial value problem (3.2) has a unique global mild solution $y(t)$, which is classical if $y_0 \in D(\mathcal{A})$. Moreover, the function $t \mapsto V(y(t))$ is non-increasing, and $\|y(t)\|_{\mathcal{H}}$ is uniformly bounded on \mathbb{R}_0^+ .*

Proof. Due to the previous results it remains to show the result for $y_0 \notin D(\mathcal{A})$. For an approximating sequence of classical solutions $\{y_n\}_{n \in \mathbb{N}}$ as in Proposition 3.6, it follows that

$$V(y(t)) = \lim_{n \rightarrow \infty} V(y_n(t)), \quad \forall t \in [0, T_{\max}(y_0)),$$

since V is continuous. Due to (2.6), we know for the classical solution $y_n(t)$ that $t \mapsto V(y_n(t))$ is non-increasing for each fixed $n \in \mathbb{N}$, i.e.

$$V(y_n(t_1)) \geq V(y_n(t_2)), \quad 0 \leq t_1 \leq t_2.$$

Letting $n \rightarrow \infty$ in this inequality shows that $t \mapsto V(y(t))$ is non-increasing as well on $[0, T_{\max})$. Hence no blow-up can occur, and thus $T_{\max}(y_0) = \infty$. \square

Definition 3.8. We define the following generalized time derivative of V for the mild solution $y(t)$ of (3.2) to the initial value $y_0 \in \mathcal{H}$:

$$\dot{V}(y_0) := \limsup_{t \searrow 0} \frac{V(y(t)) - V(y_0)}{t},$$

which may take the value $-\infty$.

Corollary 3.9. *The function $V : \mathcal{H} \rightarrow \mathbb{R}$ is a Lyapunov function for the initial value problem (3.2).*

Proof. Due to Lemma (3.2) V is continuous, and according to Theorem 3.7, we know that $t \mapsto V(y(t))$ is non-increasing for all $y_0 \in \mathcal{H}$, which proves the statement. \square

For every $y_0 \in \mathcal{H}$ we define $S(t)y_0 := y(t)$, for all $t \geq 0$, where $y(t)$ is the mild solution of (3.2) corresponding to the initial condition y_0 . Theorem 9.3.2 in [3] implies that the family $S \equiv \{S(t)\}_{t \geq 0}$ is a strongly continuous semigroup of nonlinear continuous operators in \mathcal{H} .

In the remaining part of the paper we investigate the asymptotic stability of the nonlinear semigroup S . As it turns out the semigroup is asymptotically stable “in most cases”, i.e. for all but countably many values of the parameter J . For these exceptional values of J , there exist non-trivial solutions which are periodic in time and do not decay. These states can be computed explicitly, see (5.23) below.

4. Precompactness of the trajectories. In this section we investigate the precompactness of the trajectories of (3.2). Thereby, for given $y_0 \in \mathcal{H}$ the corresponding trajectory $\gamma(y_0) \subset \mathcal{H}$ is defined by

$$\gamma(y_0) := \bigcup_{t \geq 0} S(t)y_0.$$

First, we prove the precompactness of the trajectories that are twice differentiable (in time), and then extend this result to all classical solutions. To this end we need the following lemma:

Lemma 4.1. *Let $y_0 \in D(\mathcal{A}^2)$ and let $y(t)$ be the corresponding solution of (3.2). Then $y \in C^2([0, \infty); \mathcal{H})$ and $y_t(t) \in D(\mathcal{A})$ for all $t > 0$.*

Proof. First, notice that if $y \in C^2([0, \infty), \mathcal{H})$ then $\tilde{y} := y_t$ would satisfy

$$\tilde{y}_t = A\tilde{y} + \begin{bmatrix} 0 \\ 0 \\ 0 \\ -k'_1(u(L))\frac{\psi}{m} - k'_2\left(\frac{\psi}{m}\right)\frac{\tilde{\psi}}{m} \end{bmatrix}. \quad (4.1)$$

However, for the moment we only know that $y \in C^1([0, \infty), \mathcal{H})$, see Lemma 3.5. Motivated by (4.1) we define the following functions for this fixed $y(t)$:

$$F(t) := -k'_1(u(t, L))\frac{\psi(t)}{m},$$

$$G(t, z) := -k'_2\left(\frac{\psi}{m}\right)\frac{\chi}{m} \equiv g(t)\chi,$$

where $z = [U, V, \zeta, \chi]^\top \in \mathcal{H}$. Since $y(t)$ is a classical solution, both $F(t)$ and $g(t)$ are continuously differentiable. Consequently, the operator $\tilde{\mathcal{N}} : \mathbb{R}_0^+ \times \mathcal{H} \rightarrow \mathcal{H}$, defined by $\tilde{\mathcal{N}}(t, z) := [0, 0, 0, F(t) + G(t, z)]^\top$, is also continuously differentiable (in time). Furthermore $\tilde{\mathcal{N}}$ is Lipschitz continuous in \mathcal{H} , uniformly in $t \in [0, T]$ for every $T > 0$. In the following we consider the (linear, non-autonomous) initial value problem

$$z_t = Az + \tilde{\mathcal{N}}(t, z), \quad (4.2a)$$

$$z(0) = z_0 \in \mathcal{H}. \quad (4.2b)$$

We can apply Theorem 6.1.2 in [21] which proves that there is a unique global mild solution $z(t)$ of (4.2) for every $z_0 \in \mathcal{H}$. Furthermore, if $z_0 \in D(\mathcal{A})$ then $z(t)$ is a classical solution.

Next we show that for the given classical solution $y(t)$ the function $y_t(t)$ is a mild solution of (4.2) for $z_0 = \mathcal{A}y_0$. Clearly, $y(t)$ satisfies the Duhamel formula (3.3), and differentiation with respect to t yields

$$y_t(t) = e^{tA}\mathcal{A}y_0 + \frac{d}{dt} \int_0^t e^{(t-s)A}\mathcal{N}y(s) ds. \quad (4.3)$$

According to the proof of Corollary 4.2.5 in [21] there holds

$$\frac{d}{dt} \int_0^t e^{(t-s)A}\mathcal{N}y(s) ds = e^{tA}\mathcal{N}y_0 + \int_0^t e^{(t-s)A} \frac{d}{ds} \mathcal{N}y(s) ds.$$

Inserting this in (4.3) proves that $y_t(t)$ fulfills the Duhamel formula for (4.2), and as a consequence $y_t(t)$ is the unique mild solution of (4.2) to the initial condition $z_0 = \mathcal{A}y_0$. But from the first part of the proof we know that this mild solution

$z(t) = y_t(t)$ is a classical solution of (4.2) if $\mathcal{A}y_0 \in D(\mathcal{A})$, i.e. $y_0 \in D(\mathcal{A}^2)$. So $y_t \in C^1(\mathbb{R}^+; \mathcal{H})$ and $y \in C^2(\mathbb{R}^+; \mathcal{H})$. \square

Remark 4.2. In the situation where the evolution equation is linear and autonomous, i.e. $\mathcal{N} = 0$ in our case, the above result is straightforward. If $y_0 \in D(\mathcal{A}^2)$, then we have according to Section II.5.a in [9] that $y(t) \in D(\mathcal{A}^2)$ for all $t \geq 0$. Therefore $y_t(t) = \mathcal{A}y(t) \in D(\mathcal{A})$, and so $y_{tt} = \mathcal{A}y_t = \mathcal{A}^2y$, and $y_{tt} \in C(\mathbb{R}^+)$. There it is crucial that the time derivative and the operator \mathcal{A} on the right hand side commute. This does not hold in the nonlinear situation any more, which complicates things. According to Section II.5.a in [9] the density of $D(\mathcal{A}^2)$ in \mathcal{H} is also immediate. In our case $D(\mathcal{A}^2)$ is a nonlinear subset of \mathcal{H} , see (B.6), so we need to check the density separately.

We even show the stronger property that $\mathcal{A}|_{D(\mathcal{A}^2)} \subset \mathcal{A}|_{D(\mathcal{A})}$ is dense in the product topology of $\mathcal{H} \times \mathcal{H}$.

Lemma 4.3. *For any $y \in D(\mathcal{A})$, there exists a sequence $\{y_n\}_{n \in \mathbb{N}}$ in $D(\mathcal{A}^2)$ such that $\lim_{n \rightarrow \infty} y_n = y$ and $\lim_{n \rightarrow \infty} \mathcal{A}y_n = \mathcal{A}y$ in \mathcal{H} .*

The proof was deferred to the Appendix B.

Lemma 4.4. *The trajectory $\gamma(y_0)$ is precompact for $y_0 \in D(\mathcal{A})$.*

Proof. We fix $y_0 \in D(\mathcal{A})$ and show that the corresponding trajectory $y(t)$ is precompact in \mathcal{H} . As seen in Lemma 3.5, the solution $y(t)$ is classical. Due to the compact embeddings $H^4(0, L) \hookrightarrow H^2(0, L) \hookrightarrow L^2(0, L)$ it is sufficient to show that

$$\sup_{t > 0} \|\mathcal{A}y(t)\|_{\mathcal{H}} < \infty.$$

Since $y_t = \mathcal{A}y$, this is equivalent to show that $\|y_t(t)\|_{\mathcal{H}}$ is uniformly bounded in $t > 0$.

Step 1: In the first part of this proof we assume that $y_0 \in D(\mathcal{A}^2)$. According to Lemma 4.1 the time derivative $y_t(t)$ of the corresponding solution is a classical solution of the system (2.4) differentiated in time once:

$$\rho u_{ttt} + \Lambda u_t^{\text{IV}} = 0, \quad (4.4a)$$

$$u_t(t, 0) = 0, \quad (4.4b)$$

$$u_t'(t, 0) = 0, \quad (4.4c)$$

$$m u_{ttt}(L) - \Lambda u_t'''(L) + k_1'(u(L))u_t(L) + k_2'(u_t(L))u_{tt}(L) = 0, \quad (4.4d)$$

$$J u_t'{}_{ttt}(L) + \Lambda u_t''(L) = 0. \quad (4.4e)$$

We now evaluate the time derivative of $V(y_t)$:

$$\frac{d}{dt} V(y_t) = \Lambda \int_0^L u_{tt}' u_t'' dx + \rho \int_0^L u_{ttt} u_{tt} dx + J u_t'{}_{ttt}(L) u_t'(L) \quad (4.5a)$$

$$+ m u_{ttt}(L) u_{tt}(L) + k_1(u_t(L)) u_{tt}(L) \quad (4.5b)$$

$$= u_{tt}(L) (m u_{ttt}(L) - \Lambda u_t'''(L) + k_1(u_t(L))) \quad (4.5c)$$

$$+ u_t'(L) (\Lambda u_t''(L) + J u_t'{}_{ttt}(L)) \quad (4.5d)$$

$$= u_{tt}(L) (k_1(u_t(L)) - k_1'(u(L))u_t(L) - k_2'(u_t(L))u_{tt}(L)), \quad (4.5e)$$

where we have performed partial integration in x twice and used (4.4b)-(4.4e). We have due to (2.2)

$$-k_2'(u_t(L))u_{tt}(L)^2 \leq 0, \quad \forall t \geq 0,$$

so after integration of (4.5) in time we obtain

$$V(y_t(t)) \leq V(y_t(0)) + \int_0^t u_{tt}(\tau, L) [k_1(u_t(\tau, L)) - k_1'(u(\tau, L))u_t(\tau, L)] d\tau. \quad (4.6)$$

The first integral on the right hand side, which is

$$\begin{aligned} \int_0^t u_{tt}(\tau, L) k_1(u_t(\tau, L)) d\tau &= \int_0^t \frac{d}{d\tau} \int_0^{u_t(\tau, L)} k_1(s) ds d\tau \\ &= \int_0^{u_t(t, L)} k_1(s) ds - \int_0^{u_t(0, L)} k_1(s) ds, \end{aligned} \quad (4.7)$$

is uniformly bounded since $u_t(t, L) = \frac{\psi(t)}{m}$ is uniformly bounded, see Theorem 3.7. For the remaining term in (4.6) we obtain

$$\begin{aligned} \int_0^t u_{tt}(\tau, L) k_1'(u(\tau, L)) u_t(\tau, L) d\tau &= \int_0^t \frac{d}{d\tau} \left(\frac{(u_t(\tau, L))^2}{2} \right) k_1'(u(\tau, L)) d\tau \\ &= \frac{u_t(t, L)^2}{2} k_1'(u(t, L)) - \frac{u_t(0, L)^2}{2} k_1'(u(0, L)) - \int_0^t \frac{u_t(\tau, L)^3}{2} k_1''(u(\tau, L)) d\tau. \end{aligned} \quad (4.8)$$

Due to the Sobolev embedding $H^2(0, L) \hookrightarrow C(0, L)$ we have the estimate $|u(t, L)| \leq C\|u\|_{H^2} \leq C\|y\|_{\mathcal{H}}$. Therefore $k_1''(u(t, L))$ is also uniformly bounded for $t \in [0, \infty)$. Together with the previously shown uniform boundedness of $u_t(t, L)$ we find that the first two terms in (4.8) are uniformly bounded, and for the remaining integral we get

$$\left| \int_0^t \frac{u_t(\tau, L)^3}{2} k_1''(u(\tau, L)) d\tau \right| \leq C \int_0^t |u_t(\tau, L)|^3 d\tau.$$

Due to (2.3), and considering that $u_t(t, L)$ is uniformly bounded for $t \in \mathbb{R}$, there exists a positive constant $C > 0$ such that $|k_2(u_t(t, L))| \geq C u_t(t, L)^2$ for all $t \geq 0$. This yields

$$\int_0^\infty |u_t(t, L)|^3 dt \leq C \int_0^\infty k_2(u_t(t, L)) u_t(t, L) dt,$$

and since $\frac{d}{dt}(V(y(t))) = -k_2(u_t(t, L))u_t(t, L)$ is integrable on $(0, \infty)$, we obtain $u_t(\cdot, L) \in L^3(\mathbb{R}^+)$.

Therefore, all terms in (4.8) are uniformly bounded. Together with the uniform boundedness of (4.7) this shows in (4.6) that $V(y_t(t)) \in L^\infty(\mathbb{R}^+)$, and therefore $t \mapsto \|y_t(t)\|_{\mathcal{H}}$ is uniformly bounded, see Lemma 3.3. Hence, $\gamma(y_0)$ is precompact. Moreover, notice that actually

$$\sup_{t \geq 0} \|y_t(t)\|_{\mathcal{H}} \leq \tilde{C}(\|y_0\|_{\mathcal{H}}, \|y_t(0)\|_{\mathcal{H}}), \quad (4.9)$$

where the constant \tilde{C} depends continuously on $\|y_0\|_{\mathcal{H}}$ and $\|y_t(0)\|_{\mathcal{H}}$.

Step 2: For the second part of the proof, we take $y_0 \in D(\mathcal{A})$. According to Lemma 4.3 there exists a sequence $\{y_{n,0}\}_{n \in \mathbb{N}} \subset D(\mathcal{A}^2)$ such that

$$\lim_{n \rightarrow \infty} y_{n,0} = y_0 \quad \text{and} \quad \lim_{n \rightarrow \infty} \mathcal{A}y_{n,0} = \mathcal{A}y_0. \quad (4.10)$$

For the approximating solutions $y_n(t) := S(t)y_{n,0}$ we have $(y_n)_t(0) = \mathcal{A}y_{n,0}$ for all $n \in \mathbb{N}$, and (4.10) thus implies

$$\lim_{n \rightarrow \infty} (y_n)_t(0) = \mathcal{A}y_0 \quad \text{in } \mathcal{H}. \quad (4.11)$$

Hence (4.10) and (4.11) imply that both $\{y_{n,0}\}_{n \in \mathbb{N}}$ and $\{(y_n)_t(0)\}_{n \in \mathbb{N}}$ are bounded in \mathcal{H} . Together with (4.9) this yields that

$$\sup_{\substack{t \geq 0 \\ n \in \mathbb{N}}} \|(y_n)_t(t)\|_{\mathcal{H}} < \infty,$$

i.e. $(y_n)_t$ is bounded in $L^\infty(\mathbb{R}^+; \mathcal{H})$. Now the Banach-Alaoglu Theorem, see Theorem I.3.15 in [22], shows that there exists a $w \in L^\infty(\mathbb{R}^+; \mathcal{H})$ and a subsequence $\{y_{n_k}\}_{k \in \mathbb{N}}$ such that

$$(y_{n_k})_t \overset{*}{\rightharpoonup} w \text{ in } L^\infty(\mathbb{R}^+; \mathcal{H}).$$

So for $z \in \mathcal{H}$ and $t \geq 0$ arbitrary we have

$$\lim_{k \rightarrow \infty} \int_0^t \langle (y_{n_k})_t(\tau), z \rangle_{\mathcal{H}} d\tau = \int_0^t \langle w(\tau), z \rangle_{\mathcal{H}} d\tau,$$

which is equivalent to

$$\lim_{k \rightarrow \infty} \langle y_{n_k}(t) - y_{n_k}(0), z \rangle_{\mathcal{H}} = \left\langle \int_0^t w(\tau) d\tau, z \right\rangle_{\mathcal{H}}.$$

Since $y_n(t)$ converges to $y(t)$ strongly in $L^\infty((0, T); \mathcal{H})$ for every $T > 0$, we conclude from the above

$$\langle y(t) - y(0), z \rangle_{\mathcal{H}} = \left\langle \int_0^t w(\tau) d\tau, z \right\rangle_{\mathcal{H}}.$$

Now, owing to $z \in \mathcal{H}$ being arbitrary, it follows that

$$y(t) - y(0) = \int_0^t w(\tau) d\tau. \quad (4.12)$$

Since $y \in C^1(\mathbb{R}^+; \mathcal{H})$, we can take the time derivative of (4.12), and obtain $y_t \equiv w$. This implies $y_t \in L^\infty(\mathbb{R}^+; \mathcal{H})$, i.e. $\|y_t(\cdot)\|_{\mathcal{H}}$ is uniformly bounded, which proves the precompactness of $\gamma(y_0)$. \square

Remark 4.5. In the linear case, i.e. $\mathcal{N} = 0$, the proof of the trajectory precompactness is much simpler: For classical solutions $y(t)$ we have $Ay(t) = Ae^{tA}y_0 = e^{tA}Ay_0$, so $Ay(t)$ is uniformly bounded. Since A^{-1} is compact, this proves the precompactness for classical solutions. Since e^{tA} is a contraction semigroup, any mild solution can be approximated uniformly by classical solutions, and the precompactness follows also for mild solutions.

In the case when k_1 is linear, and k_2 is nonlinear, the precompactness property of the trajectories can also be verified easily: If we incorporate the k_1 -term and the linear part of k_2 into A , this operator still generates a contraction semigroup, it is invertible and has a compact resolvent. For the remaining nonlinear term we can show $\mathcal{N}(y(\cdot)) \in L^1(\mathbb{R}^+; \mathcal{H})$, using (3.4) and (5.1). Then the prerequisites of Theorem 4 in the article by Dafermos and Slemrod [8] are fulfilled, and the precompactness of the trajectories for all mild solutions follows.

5. ω -limit set and asymptotic stability. In this section we first investigate some properties of ω -limit sets, which will turn out to be essential in the study of the asymptotic stability of the system.

Definition 5.1. Given the semigroup S , the ω -limit set for $y_0 \in \mathcal{H}$ is denoted by $\omega(y_0)$, and defined by:

$$\omega(y_0) := \{y \in \mathcal{H} : \exists \{t_n\}_{n \in \mathbb{N}} \subset \mathbb{R}^+, \lim_{n \rightarrow \infty} t_n = \infty \wedge \lim_{n \rightarrow \infty} S(t_n)y_0 = y\}$$

It is possible that $\omega(y_0) = \emptyset$. According to Proposition 9.1.7 in [3] we have:

Lemma 5.2. *For $y_0 \in \mathcal{H}$, the set $\omega(y_0)$ is S -invariant, i.e. $S(t)\omega(y_0) \subset \omega(y_0)$ for all $t \geq 0$.*

According to (3.9) the function $t \mapsto V(S(t)y_0)$ is monotonically (but not necessarily strictly) decreasing for any fixed $y_0 \in \mathcal{H}$. Furthermore it is bounded from below by 0. Therefore, the following limit exists:

$$\nu(y_0) := \lim_{t \rightarrow \infty} V(S(t)y_0) \geq 0. \quad (5.1)$$

Lemma 5.3. *Suppose $\omega(y_0) \neq \emptyset$. Then there holds*

$$V(\omega(y_0)) = \{\nu(y_0)\},$$

i.e. V takes the same value $\nu(y_0)$ on every element of $\omega(y_0)$. In particular $\dot{V}(y) = 0$ for all $y \in \omega(y_0)$.

Proof. Let $y \in \omega(y_0)$. There exists a sequence $\{t_n\}_{n \in \mathbb{N}} \subset \mathbb{R}^+$, with $t_n \rightarrow \infty$, such that $\lim_{n \rightarrow \infty} S(t_n)y_0 = y$. Since V is continuous,

$$\lim_{n \rightarrow \infty} V(S(t_n)y_0) = V(y).$$

Due to (5.1) we have $V(y) = \nu(y_0)$, and the result follows. \square

Therefore we may identify possible ω -limit sets by investigating trajectories along which the Lyapunov function V is constant. To this end we define the set $\Omega \subset \mathcal{H}$ as the largest S -invariant subset of $\{y \in \mathcal{H} : \dot{V}(y) = 0\}$. Clearly there holds

$$\omega(y_0) \subset \Omega, \quad \forall y_0 \in \mathcal{H}.$$

Thus our focus first lies on characterizing Ω .

Lemma 5.4. *For every $y_0 \in \mathcal{H}$ the following holds, for all $t > 0$:*

$$\int_0^t S(s)y_0 \, ds \in D(\mathcal{A}), \quad (5.2)$$

and

$$S(t)y_0 - y_0 = A \int_0^t S(s)y_0 \, ds + \int_0^t \mathcal{N}S(s)y_0 \, ds. \quad (5.3)$$

The proof was deferred to the Appendix B. This result can be understood as a generalization of Theorem 1.2.4 in [21] to nonlinear semigroups. It even holds in the general situation where A is linear and the infinitesimal generator of a C_0 -semigroup, and \mathcal{N} is differentiable, see [23] for a more details.

Proposition 5.5. *For all $y = [u, v, \xi, \psi]^\top \in \Omega$ there holds $\psi = 0$, $u(L) = 0$.*

Proof. For a fixed $y_0 \in \Omega$ let $y(t) = S(t)y_0$. Since Ω is S -invariant we have $V(y(t)) = \nu(y_0)$ for all $t \geq 0$. First we show that

$$\psi(t) = 0, \quad \forall t \geq 0. \quad (5.4)$$

In the case when $y_0 \in \Omega \cap D(\mathcal{A})$, (5.4) follows easily since (2.6) implies for the corresponding classical solution

$$\dot{V}(y(t)) = 0 \quad \Leftrightarrow \quad \psi(t) = 0.$$

Next we investigate the case when $y_0 \in \Omega \setminus D(\mathcal{A})$. Then there is a sequence $\{y_{n,0}\}_{n \in \mathbb{N}} \subset D(\mathcal{A})$ such that $\lim_{n \rightarrow \infty} y_{n,0} = y_0$ in \mathcal{H} . Theorem 3.7 implies $y_n(t) \rightarrow$

$y(t)$ in $C([0, T]; \mathcal{H})$ for any $T > 0$, where $y_n(t) = S(t)y_{n,0}$. Since V is locally Lipschitz continuous in \mathcal{H} , $\{V(y_n(t))\}_{n \in \mathbb{N}}$ is a Cauchy sequence in $C([0, T]; \mathbb{R})$. On the other hand we also have the convergence

$$\psi_n(t) \rightarrow \psi(t) \quad \text{in } C([0, T]; \mathbb{R}). \quad (5.5)$$

Due to (2.6) this implies that

$$\left\{ \frac{d}{dt} V(y_n(t)) \right\}_{n \in \mathbb{N}}$$

is a Cauchy sequence in $C([0, T]; \mathbb{R})$. We conclude that $\{V(y_n(t))\}_{n \in \mathbb{N}}$ is a Cauchy sequence in $C^1([0, T]; \mathbb{R})$. So there exists a unique $w \in C^1([0, T]; \mathbb{R})$ such that

$$V(y_n(t)) \rightarrow w(t) \quad \text{in } C^1([0, T]; \mathbb{R}). \quad (5.6)$$

On the other hand, we know that $V(y_n(t)) \rightarrow V(y(t)) = \nu(y_0)$ for every $t \geq 0$, and hence $w(t) \equiv \nu(y_0)$. Together with (5.6) this implies $\dot{V}(y_n(t)) = -k_2 \left(\frac{\psi_n}{m}\right) \frac{\psi_n}{m} \rightarrow 0$ uniformly on $[0, T]$. With (5.5) this now yields (5.4) and in particular $\psi(0) = 0$.

Next, $u(t, L) = 0$ for all $t \geq 0$ is demonstrated. From (5.2) and (5.4) it follows that

$$m \left(\int_0^t v(s) \, ds \right) \Big|_{x=L} = \int_0^t \psi(s) \, ds = 0.$$

Using this, the first component of (5.3) implies

$$0 = \left(\int_0^t v(s) \, ds \right) \Big|_{x=L} = u(t, L) - u(0, L).$$

Therefore $u(t, L)$ is constant along $y(t)$, which implies

$$\int_0^t u(s, L) \, ds = u_0(L)t, \quad t \geq 0. \quad (5.7)$$

Since $\sup_{t>0} \|y(t)\|_{\mathcal{H}} < \infty$, it follows that $\sup_{t>0} \|v(t)\|_{L^2(0,L)} < \infty$. Therefore, the second component of (5.3) implies

$$\sup_{t \geq 0} \left\| \left(\int_0^t u(s) \, ds \right)^{\text{IV}} \right\|_{L^2(0,L)} < \infty. \quad (5.8)$$

Next we apply the following Gagliardo-Nirenberg inequality (cf. [18]), which guarantees the existence of $C > 0$ such that there holds for all $t \geq 0$:

$$\left\| \int_0^t u(s) \, ds \right\|_{L^\infty(0,L)} \leq C \left\| \left(\int_0^t u(s) \, ds \right)^{\text{IV}} \right\|_{L^2(0,L)}^{\frac{1}{8}} \left\| \int_0^t u(s) \, ds \right\|_{L^2(0,L)}^{\frac{7}{8}}. \quad (5.9)$$

The first factor on the right hand side is uniformly bounded due to (5.8). For the second factor we observe that, according to Theorem 3.7, $t \mapsto \|u(t)\|_{L^2(0,L)}$ is uniformly bounded, and therefore $t \mapsto \left\| \int_0^t u(s) \, ds \right\|_{L^2(0,L)}$ grows at most linearly. Altogether this implies in (5.9) that $t \mapsto \int_0^t u(s, L) \, ds$ grows at most like $t^{\frac{7}{8}}$. But this contradicts (5.7) unless $u_0(L) = 0$. This shows that $u(t, L) = 0$ for all $t \geq 0$. \square

This result allows to represent any trajectory $\gamma(y_0) \subset \Omega$ as a solution to a simpler linear system characterizing Ω . By inserting the result of Proposition 5.5

in the equation (5.3) we find that any mild solution $y(t)$ of (3.2) with $y(t) \in \Omega$ for all $t \geq 0$, satisfies the following system:

$$u(t) - u(0) = \int_0^t v(s) \, ds, \quad (5.10a)$$

$$v(t) - v(0) = -\frac{\Lambda}{\rho} \left(\int_0^t u(s) \, ds \right)^{\text{IV}}, \quad (5.10b)$$

$$\xi(t) - \xi(0) = -\Lambda \left(\int_0^t u(s) \, ds \right)'' \Big|_{x=L}, \quad (5.10c)$$

$$0 = \left(\int_0^t u(s) \, ds \right)''' \Big|_{x=L}, \quad (5.10d)$$

together with the additional boundary condition $u(t, L) = 0$. We will show that this system is overdetermined. To this end we first investigate the system (5.10) without the condition $u(t, L) = 0$, and only incorporate it later.

The system (5.10a)-(5.10c) can be interpreted as a mild formulation of a linear evolution equation in a Hilbert space $\tilde{\mathcal{H}}$:

$$w_t = \mathcal{B}w, \quad (5.11)$$

with $w = [u, v, \xi]^\top \in \tilde{\mathcal{H}}$. Thereby $\tilde{\mathcal{H}}$ is the Hilbert space

$$\tilde{\mathcal{H}} := \{w = [u, v, \xi]^\top : u \in \tilde{H}_{0,\mathbb{R}}^2(0, L), v \in L_{\mathbb{R}}^2(0, L), \xi \in \mathbb{R}\},$$

and \mathcal{B} is the following linear operator in $\tilde{\mathcal{H}}$:

$$\mathcal{B} \begin{bmatrix} u \\ v \\ \xi \end{bmatrix} = \begin{bmatrix} v \\ -\frac{\Lambda}{\rho} u^{\text{IV}} \\ -\Lambda u''(L) \end{bmatrix}, \quad (5.12)$$

with the domain

$$D(\mathcal{B}) := \{w \in \tilde{\mathcal{H}} : u \in \tilde{H}_{0,\mathbb{R}}^4(0, L), v \in \tilde{H}_{0,\mathbb{R}}^2(0, L), \xi = Jv'(L), u'''(L) = 0\},$$

which incorporates the condition (5.10d). The space $\tilde{\mathcal{H}}$ is equipped with the inner product

$$\langle\langle w_1, w_2 \rangle\rangle := \frac{\Lambda}{2} \int_0^L u_1'' u_2'' \, dx + \frac{\rho}{2} \int_0^L v_1 v_2 \, dx + \frac{1}{2J} \xi_1 \xi_2.$$

Due to Proposition A.3 the operator \mathcal{B} is skew-adjoint (in $\tilde{\mathcal{X}}$, i.e. the complexification of $\tilde{\mathcal{H}}$, see the Appendix A) and generates a C_0 -group of unitary operators. The eigenvalues $\{\mu_n\}_{n \in \mathbb{Z} \setminus \{0\}}$ are purely imaginary, and come in complex conjugated pairs, i.e. $\mu_{-n} = \overline{\mu_n}$. Zero is no eigenvalue, since \mathcal{B} is invertible, see [14]. The corresponding eigenfunctions $\{\Phi_n\}_{n \in \mathbb{Z} \setminus \{0\}}$ form an orthonormal basis of $\tilde{\mathcal{X}}$. They are given by

$$\Phi_n = \begin{bmatrix} u_n \\ \mu_n u_n \\ \mu_n J u_n'(L) \end{bmatrix}, \quad (5.13)$$

where u_n is the unique real-valued solution of

$$\rho \mu_n^2 u_n + \Lambda u_n^{\text{IV}} = 0, \quad (5.14a)$$

$$u_n'''(L) = 0, \quad (5.14b)$$

$$J \mu_n^2 u_n'(L) + \Lambda u_n''(L) = 0, \quad (5.14c)$$

where u_n is normalized such that $\|\Phi_n\|_{\tilde{\mathcal{X}}} = 1$. Note that $\mu_n^2 < 0$. From (5.13) it is clear that $\Phi_{-n} = \overline{\Phi_n}$, and hence $u_{-n} = u_n$. For the complete spectral analysis of \mathcal{B} see Proposition A.3 in the appendix. For notational simplicity we include the index $n = 0$ in the following by setting $\mu_0 := 0$ and $\Phi_0 := 0$ and $u_0 := 0$.

For our further considerations we need the following lemma.

Lemma 5.6. *There exists a non-trivial solution u_n of the system (5.14) that additionally satisfies $u_n(L) = 0$ iff*

$$J = \rho \left(\frac{L}{\ell\pi} \right)^3 \frac{(-1)^\ell + \cosh \ell\pi}{\sinh \ell\pi}, \quad \text{for some } \ell \in \mathbb{N}.$$

In this case, $u_n (= u_{-n})$ is unique up to normalization and $\mu_n^2 = -\frac{\Lambda}{\rho} \left(\frac{\ell\pi}{L} \right)^4$. We shall denote the index of this particular eigenfunction by $n = n^*(\ell) > 0$.

The proof is deferred to the Appendix B. For further use we define the set

$$\mathcal{J} := \left\{ \rho \left(\frac{L}{\ell\pi} \right)^3 \frac{(-1)^\ell + \cosh \ell\pi}{\sinh \ell\pi} : \ell \in \mathbb{N} \right\}.$$

We denote the ℓ -th entry by J_ℓ .

Theorem 5.7. *Concerning the ω -limit set we distinguish between two situations:*

- (i) *Assume that the given model parameter $J \notin \mathcal{J}$. Then $w = [u, v, \xi]^\top \equiv 0$ is the only solution to (5.11) with $u(L) = 0$, and therefore $\Omega = \{0\}$.*
- (ii) *If $J \in \mathcal{J}$, then Ω is*

$$\text{span}_{\mathbb{R}} \{ [u_{n^*}, 0, 0, 0]^\top, [0, u_{n^*}, Ju'_{n^*}(L), 0]^\top \}.$$

Thereby u_{n^} is the non-trivial solution from Lemma 5.6.*

Proof. This proof closely follows the argumentation in [5]. According to Proposition A.3 in the Appendix A we can write the mild solution of the linear evolution equation (5.11) with the initial condition $w_0 \in \tilde{\mathcal{H}}$ as

$$w(t) = e^{t\mathcal{B}} w_0 = \sum_{n \in \mathbb{Z}} \langle w_0, \Phi_n \rangle_{\tilde{\mathcal{X}}} e^{\mu_n t} \Phi_n, \quad (5.15)$$

where $\{\mu_n\}_{n \in \mathbb{Z}}$ are the (imaginary) eigenvalues of \mathcal{B} , and the Φ_n are the corresponding normalized eigenfunctions¹, see Proposition A.3. Thereby $\langle \cdot, \cdot \rangle_{\tilde{\mathcal{X}}}$ is the inner product in $\tilde{\mathcal{X}}$, see the Appendix A. We define $c_n := \langle w_0, \Phi_n \rangle_{\tilde{\mathcal{X}}}$ for all $n \in \mathbb{Z}$. Due to the orthonormality of the eigenfunctions $\{\Phi_n\}_{n \in \mathbb{Z} \setminus \{0\}}$ and the fact that $\{\mu_n\}_{n \in \mathbb{Z}} \subset i\mathbb{R}$ we have for any $N \in \mathbb{N}$:

$$\left\| \sum_{|n| \geq N} c_n e^{\mu_n t} \Phi_n \right\|_{\tilde{\mathcal{X}}}^2 = \sum_{|n| \geq N} |c_n|^2. \quad (5.16)$$

Due to Parseval's identity we also have $\sum_{n \in \mathbb{Z}} |\langle w_0, \Phi_n \rangle_{\tilde{\mathcal{X}}}|^2 = \|w_0\|_{\tilde{\mathcal{X}}}^2$. As a consequence the right hand side in (5.16) tends to zero as $N \rightarrow \infty$. So, for every $\varepsilon > 0$ there exists some $N > 0$ such that

$$\sup_{t \geq 0} \left\| \sum_{|n| \geq N} c_n e^{\mu_n t} \Phi_n \right\|_{\tilde{\mathcal{X}}} < \varepsilon. \quad (5.17)$$

¹Note that if $w_0 \in \tilde{\mathcal{H}}$, i.e. w_0 is real valued, then the series always maps into $\tilde{\mathcal{H}}$ again.

The first component of the series (5.15) converges in $H^2(0, L)$ and therefore also in $C([0, L])$. Thus we have

$$u(t, L) = \sum_{n \in \mathbb{Z}} c_n e^{\mu_n t} u_n(L), \quad \forall t \geq 0. \quad (5.18)$$

Using this representation formula we now investigate those $u(t)$ that satisfy $u(t, L) = 0$ for all times. We immediately find for every $N \in \mathbb{N}$:

$$\begin{aligned} \left| \sum_{|n| \geq N} c_n e^{\mu_n t} u_n(L) \right| &\leq C \left\| \sum_{|n| \geq N} c_n e^{\mu_n t} u_n \right\|_{H^2(0, L)} \\ &\leq C \left\| \sum_{|n| \geq N} c_n e^{\mu_n t} \Phi_n \right\|_{\tilde{\mathcal{X}}}. \end{aligned}$$

According to (5.17) this implies that, for every $\varepsilon > 0$, we can find an $N \in \mathbb{N}$ (large enough) such that

$$\sup_{t \geq 0} \left| \sum_{n=-N}^N c_n e^{\mu_n t} u_n(L) \right| < \varepsilon, \quad (5.19)$$

provided that $u(t, L) = 0$ for all $t \geq 0$.

We fix now some $k \in \mathbb{Z}$ and $\varepsilon > 0$, and select $N \in \mathbb{N}$ so large that $|k| < N$ and (5.19) is satisfied. Then we multiply the finite sum by $e^{-\mu_k t}$ and integrate over $[0, T]$:

$$\frac{1}{T} \int_0^T \sum_{n=-N}^N c_n e^{\mu_n t} u_n(L) e^{-\mu_k t} dt = \sum_{n=-N}^N c_n u_n(L) \frac{1}{T} \int_0^T e^{(\mu_n - \mu_k) t} dt.$$

Due to (5.19) this expression still has modulus less than ε . Now we let $T \rightarrow \infty$. Since all eigenvalues μ_n of \mathcal{B} are distinct (see Proposition A.3), all terms in the integral vanish except for the term where $n = k$, and we obtain

$$|c_k u_k(L)| < \varepsilon.$$

Since ε was arbitrary, we conclude

$$c_k u_k(L) = 0, \quad \forall k \in \mathbb{Z}. \quad (5.20)$$

Now we need to distinguish between two situations: Either $J \notin \mathcal{J}$ or $J \in \mathcal{J}$.

(i) In the first case, due to Lemma 5.6, $u_n(L) \neq 0$ for all $n \in \mathbb{Z}$. Then (5.20) implies that $c_k = 0$ for all $k \in \mathbb{Z}$, and consequently $w_0 = w(t) \equiv 0$ for all $t > 0$. Therefore $\Omega = \{0\}$.

(ii) Now we consider $J = J_\ell \in \mathcal{J}$. According to Lemma 5.6 we have $u_k(L) = 0$ iff $k \neq \pm n^*(\ell)$. So we get from (5.20) that

$$c_k = 0, \quad \forall k \in \mathbb{Z} \setminus \{\pm n^*(\ell)\}, \quad (5.21a)$$

$$c_{n^*} \in \mathbb{C} \text{ arbitrary}, \quad (5.21b)$$

and $c_{-n^*} = \overline{c_{n^*}}$. Together with $\psi = 0$ in Ω we find in this case that $\Omega = \text{Respan}_{\mathbb{C}}\{[\Phi_{-n^*}, 0]^\top, [\Phi_{n^*}, 0]^\top\} = \text{span}_{\mathbb{R}}\{[u_{n^*}, 0, 0, 0]^\top, [0, u_{n^*}, Ju'_{n^*}(L), 0]^\top\}$. \square

Remark 5.8. An alternative approach is to consider the system (5.10a)-(5.10c) together with $u(t, L) = 0$, ignoring (5.10d) for the moment. The system (5.11) is then defined in

$$\tilde{\mathcal{H}}_1 := \{w \in \tilde{\mathcal{H}} : u(L) = 0\}$$

instead of $\tilde{\mathcal{H}}$, and \mathcal{B} has a different domain:

$$D_1(\mathcal{B}) := \{w \in \tilde{\mathcal{H}}_1 : u \in \tilde{H}_{0,\mathbb{R}}^4(0, L), v \in \tilde{H}_{0,\mathbb{R}}^2(0, L), \xi = Jv'(L), v(L) = 0\}.$$

Analogously to the Proposition A.3 one finds that the operator $(\mathcal{B}, D_1(\mathcal{B}))$ is again skew-adjoint, generates a C_0 -semigroup of unitary operators, and its eigenfunctions form an orthogonal basis. For the first component u_n of the eigenfunctions we again get the representation (A.1). However, here we use $u_n(L) = 0$ in order to determine the constants (i.e. the u_n are in general different to the ones used in the proof above). With these u_n we have again (5.18), and only there we apply the remaining condition (5.10d).

In the case $J \in \mathcal{J}$ we have seen (in Theorem 5.7) that $\Omega = \text{Re span}\{[\Phi_{\pm n^*}, 0]^\top\}$. From the definition of the $\Phi_{\pm n^*}$ we find that they are precisely the (two) common eigenfunctions of $(\mathcal{B}, D(\mathcal{B}))$ and $(\mathcal{B}, D_1(\mathcal{B}))$. We conclude that, in order to determine the ω -limit set, the two approaches using either $(\mathcal{B}, D(\mathcal{B}))$ or $(\mathcal{B}, D_1(\mathcal{B}))$ are equivalent. They only differ in the order in which the boundary conditions $u'''(L) = 0$ and $u(L) = 0$ are applied.

Now we have all the prerequisites to prove our main result.

Theorem 5.9. *Assume $J \notin \mathcal{J}$. For every $y_0 \in D(\mathcal{A})$,*

$$\lim_{t \rightarrow \infty} y(t) = 0,$$

i.e. the system (3.2) is asymptotically stable with respect to $\|\cdot\|_{\mathcal{H}}$.

Proof. Due to Lemma 4.4, the trajectory $\gamma(y_0)$ is precompact. According to Theorem 5.7 we further have $\Omega = \{0\}$. So we can apply the LaSalle Invariance Principle, cf. Theorem 3.64 in [16], which proves that $\lim_{t \rightarrow \infty} \|y(t)\|_{\mathcal{H}} = 0$. \square

Theorem 5.10. *Let $J = J_\ell \in \mathcal{J}$ for some $\ell \in \mathbb{N}$. Given an initial condition $y_0 \in D(\mathcal{A})$ the corresponding solution $y(t)$ of (3.2) approaches (with respect to $\|\cdot\|_{\mathcal{H}}$) the time-periodic solution corresponding to the initial condition $\Pi^* y_0$ as $t \rightarrow \infty$. Thereby Π^* is the orthogonal projection from \mathcal{H} onto Ω , and it is given by*

$$\Pi^* y = \begin{bmatrix} \Lambda \langle u'', u_{n^*}'' \rangle_{L^2} u_{n^*} \\ |\mu_{n^*}|^2 (\rho \langle v, u_{n^*} \rangle_{L^2} + \xi u_{n^*}'(L)) u_{n^*} \\ J |\mu_{n^*}|^2 (\rho \langle v, u_{n^*} \rangle_{L^2} + \xi u_{n^*}'(L)) u_{n^*}'(L) \\ 0 \end{bmatrix}, \quad (5.22)$$

where $\langle \cdot, \cdot \rangle_{L^2}$ denotes the standard inner product on $L^2(0, L)$.

Proof. Let $n^*(\ell)$ be as in Lemma 5.6. According to the end of the proof of Theorem 5.7 the ω -limit set is the (complex) span of the two vectors $\Psi_{\pm n^*} = [\Phi_{\pm n^*}, 0]^\top$, where $\Phi_{-n^*} = \overline{\Phi_{n^*}}$. Since $\Phi_{\pm n^*} \in D(\mathcal{B})$ we know that $u_{n^*}'''(L) = 0$, and so the $\Psi_{\pm n^*}$ are eigenvectors of A to the eigenvalues $\pm \mu_{n^*}$. We may now define the orthogonal projection (first in \mathcal{X} , see the Appendix A):

$$\Pi^* := \langle \cdot, \Psi_{-n^*} \rangle_{\mathcal{X}} \Psi_{-n^*} + \langle \cdot, \Psi_{n^*} \rangle_{\mathcal{X}} \Psi_{n^*}.$$

According to Proposition A.1 the eigenvectors of A form an orthogonal basis of \mathcal{X} , so Π^* commutes with A , and $\mathcal{X} = \ker \Pi^* \oplus \text{ran } \Pi^*$ is an orthogonal, A -invariant decomposition of \mathcal{X} . In the following we work with the restriction of Π^* to \mathcal{H} , and keep the same notation. The explicit representation of Π^* is given by (5.22).

In the next step we show that Π^* commutes with the nonlinearity \mathcal{N} . Since the first component u_{n^*} of Ψ_{n^*} satisfies $u_{n^*}(L) = 0$, it is clear that $\mathcal{N}\Psi_{\pm n^*} = 0$ and thus $\mathcal{N}\Pi^* = 0$. Let now $y \in \mathcal{X}$. Then

$$\mathcal{N}y = \begin{bmatrix} 0 \\ 0 \\ 0 \\ -k_1(u(L)) - k_2\left(\frac{\psi}{m}\right) \end{bmatrix},$$

and so $\Pi^*\mathcal{N}y = 0$.

As a consequence, the decomposition $\mathcal{H} = \ker \Pi^* \oplus \text{ran } \Pi^*$ is invariant under the nonlinear semigroup S generated by \mathcal{A} . The trajectories of $S|_{\ker \Pi^*}$ lying in $D(\mathcal{A})$ are still precompact. We know from Theorem 5.7 that any ω -limit set of $S|_{\ker \Pi^*} \subset S$ has to be a subset of $\text{ran } \Pi^*$. But on the other hand any trajectory and limit of $S|_{\ker \Pi^*}$ has to lie within $\ker \Pi^*$, which is orthogonal to $\text{ran } \Pi^*$. Thus the only possible ω -limit set for $S|_{\ker \Pi^*}$ is $\{0\} = \text{ran } \Pi^* \cap \ker \Pi^*$. And therefore $S(t)y_0$ approaches $S(t)\Pi^*y_0$ as $t \rightarrow \infty$. \square

Remark 5.11. The asymptotic limit described in Theorem 5.10 can be computed explicitly. If $J = J_\ell$ for some $\ell \in \mathbb{N}$, it follows from (5.18), (5.21) and Lemma 5.6 that all real non-decaying solutions u_p of (2.4) are given by

$$u_p(t, x) = T(t)u_{n^*}(x), \quad (5.23)$$

where

$$T(t) = a \cos \sqrt{\frac{\Lambda}{\rho}} \left(\frac{\ell\pi}{L} \right)^2 t + b \sin \sqrt{\frac{\Lambda}{\rho}} \left(\frac{\ell\pi}{L} \right)^2 t, \quad a, b \in \mathbb{R},$$

and u_{n^*} is given by (B.16). In particular, it follows from Theorem 5.10 that for a given initial condition y_0 the solution u of (2.4) approaches the solution u_p given in (5.23), with the coefficients a and b determined by:

$$a := \Lambda \langle u_0'', u_{n^*}'' \rangle_{L^2},$$

and

$$b := -\sqrt{\frac{\Lambda}{\rho}} \left(\frac{\ell\pi}{L} \right)^2 (\rho \langle v_0, u_{n^*} \rangle_{L^2} + \xi_0 u_{n^*}'(L)).$$

Remark 5.12. As already mentioned in Remark 4.5 it is the nonlinear term $k_1(u(L))$, representing the spring in the model, that prevents the nonlinear operator \mathcal{A} from being dissipative. As a consequence the semigroup S is not contractive and therefore it is not possible to extend the precompactness of the classical trajectories to the trajectories of the mild solutions using a density argument (for this see the proof of Theorem 3.65 in [16]). From the physical point of view one might expect that (at least for $J \notin \mathcal{J}$) also the mild solutions tend to zero, which is motivated by the observation that the total energy is dissipated whenever a trajectory does not lie in $\Omega = \{0\}$, i.e. for almost all times the system loses energy due to friction. However, from the mathematical point of view it is not clear that the trajectory converges at all as $t \rightarrow \infty$. This is due to the missing information if the trajectory is precompact for a non-classical solution.

Appendix A. Functional analytical results. Even though the analysis of this paper is carried out for real-valued functions u and as a consequence in the real Hilbert space \mathcal{H} , the spectral analysis of the occurring linear operators needs to be performed in a complex Hilbert space. This section contains some of those results. For the spectral analysis of the operator A we introduce the complex Hilbert space

$$\mathcal{X} := \{y = [u, v, \xi, \psi]^\top : u \in \tilde{H}_0^2(0, L), v \in L^2(0, L), \xi, \psi \in \mathbb{C}\},$$

equipped with the inner product

$$\langle y_1, y_2 \rangle_{\mathcal{X}} := \frac{\Lambda}{2} \int_0^L u_1'' \overline{u_2''} dx + \frac{\rho}{2} \int_0^L v_1 \overline{v_2} dx + \frac{1}{2J} \xi_1 \overline{\xi_2} + \frac{1}{2m} \psi_1 \overline{\psi_2}, \quad \forall y_1, y_2 \in \mathcal{X}.$$

For the operator A given by (3.1) we consider the natural continuation to \mathcal{X} , still denoted by A . This continuation still is of the form (3.1), and the domain is now

$$D_{\mathbb{C}}(A) = \{y \in \mathcal{X} : u \in \tilde{H}_0^4(0, L), v \in \tilde{H}_0^2(0, L), \xi = Jv'(L), \psi = mv(L)\},$$

where the occurring Sobolev spaces now consist of complex valued functions.

Proposition A.1. *The linear operator A is skew-adjoint and has compact resolvent in \mathcal{X} . The spectrum $\sigma(A)$ consists of countably many eigenvalues $\{\lambda_n\}_{n \in \mathbb{Z}}$. They are all isolated and purely imaginary, and each eigenspace has finite dimension. The eigenspaces form a complete orthogonal decomposition of \mathcal{X} .*

Proof. It can easily be shown that for all $y_1, y_2 \in D_{\mathbb{C}}(A)$

$$\langle Ay_1, y_2 \rangle_{\mathcal{X}} = \frac{\Lambda}{2} \int_0^L v_1'' \overline{u_2''} - u_1'' \overline{v_2''} dx = -\langle y_1, Ay_2 \rangle_{\mathcal{X}},$$

i.e. A is skew-symmetric. Straightforward calculations, analogous to those in [14], demonstrate that A is invertible and $A^{-1} : \mathcal{X} \rightarrow \mathcal{X}$ is even compact. So $0 \in \rho(A)$, and due to the corollary of Theorem VII.3.1 in [26] this proves that A is skew-adjoint. Then, according to Theorem III.6.26 in [13] the spectrum $\sigma(A)$ consists of countably many eigenvalues, which are all isolated. The corresponding eigenspaces are finite-dimensional, and the eigenvectors form an orthogonal basis according to Theorem V.2.10 in [13]. \square

The following is an extension of Lemma 3.1 from \mathcal{H} to \mathcal{X} .

Lemma A.2. *The linear operator A generates a C_0 -semigroup $(e^{tA})_{t \geq 0}$ of unitary operators in \mathcal{X} .*

Proof. From Proposition A.1 we know that A is skew-adjoint in \mathcal{X} . So we may apply Stone's Theorem, and $(e^{tA})_{t \geq 0}$ is a C_0 -semigroup of unitary operators in \mathcal{X} . \square

Next we turn to the spectral analysis of \mathcal{B} . To this end we introduce the Hilbert space

$$\tilde{\mathcal{X}} := \{w = [u, v, \xi]^\top : u \in \tilde{H}_0^2(0, L), v \in L^2(0, L), \xi \in \mathbb{C}\},$$

equipped with the inner product

$$\langle\langle w_1, w_2 \rangle\rangle_{\tilde{\mathcal{X}}} := \frac{\Lambda}{2} \int_0^L u_1'' \overline{u_2''} dx + \frac{\rho}{2} \int_0^L v_1 \overline{v_2} dx + \frac{1}{2J} \xi_1 \overline{\xi_2}.$$

The continuation of \mathcal{B} to $\tilde{\mathcal{X}}$ is still denoted by \mathcal{B} and given by (5.12), and has the domain

$$D_{\mathbb{C}}(\mathcal{B}) := \{y \in \tilde{\mathcal{X}} : u \in \tilde{H}_0^4(0, L), v \in \tilde{H}_0^2(0, L), \xi = Jv'(L), u'''(L) = 0\}.$$

Proposition A.3. *The operator \mathcal{B} is skew-adjoint and has compact resolvent in $\tilde{\mathcal{X}}$. The spectrum $\sigma(\mathcal{B})$ consists entirely of isolated eigenvalues $\{\mu_n\}_{n \in \mathbb{Z} \setminus \{0\}}$ located on the imaginary axis, and they have no accumulation point. All eigenspaces are one-dimensional, and the corresponding eigenfunctions form an orthogonal basis of $\tilde{\mathcal{X}}$. The normalized eigenfunction associated to μ_n is given by*

$$\Phi_n = \begin{bmatrix} u_n \\ \mu_n u_n \\ \mu_n J u_n'(L) \end{bmatrix}, \quad n \in \mathbb{Z} \setminus \{0\},$$

where the real function $u_n \in \tilde{H}_0^4(0, L)$ is the unique (up to normalization) solution of the boundary value problem (5.14). Thereby u_n is scaled such that $\|\Phi_n\|_{\tilde{\mathcal{X}}} = 1$.

Proof. Analogously to the proof of Proposition A.1 we show that $0 \in \rho(\mathcal{B})$, that \mathcal{B}^{-1} is compact in $\tilde{\mathcal{X}}$ and that \mathcal{B} is skew-adjoint. Now we can apply the Corollary of Theorem VII.3.1 in [26], which proves that the skew-symmetric operator \mathcal{B} is even skew-adjoint. According to Theorem III.6.26 in [13] the spectrum $\sigma(\mathcal{B})$ consists of countably many eigenvalues $\{\mu_n\}_{n \in \mathbb{Z} \setminus \{0\}}$, which are all isolated. They come in complex conjugated pairs, i.e. $\mu_{-n} = \overline{\mu_n}$. The corresponding eigenspaces are finite-dimensional, and the eigenvectors form an orthogonal basis according to Theorem V.2.10 in [13]. Since \mathcal{B} is skew-adjoint we have $\sigma(\mathcal{B}) \subset i\mathbb{R}$. Finally, the fact that the $\{\Phi_n\}_{n \in \mathbb{Z} \setminus \{0\}}$ are an orthogonal basis of $\tilde{\mathcal{X}}$ follows immediately from the application of Theorem V.2.10 in [13].

Let $\Phi_n = [u_n, v_n, \xi_n]^\top \in D_{\mathbb{C}}(\mathcal{B})$ be an eigenfunction corresponding to μ_n for $n \in \mathbb{Z} \setminus \{0\}$, i.e. $\mathcal{B}\Phi_n = \mu_n \Phi_n$. Now Φ_n satisfies the eigenvalue equation iff u_n solves (5.14). The v_n and ξ_n can be determined from u_n via $v_n = \mu_n u_n$ and $\xi_n = J \mu_n u_n'(L)$. The system (5.14) has a non-trivial solution iff $\mu_n \in \sigma(\mathcal{B})$ (note that we have already shown that $0 \notin \sigma(\mathcal{B})$, i.e. we may assume $\mu_n \neq 0$). In this case we get the general solution $u_n \in \tilde{H}_0^4(0, L)$ of (5.14a) as

$$u_n(x) = C_1[\cosh px - \cos px] + C_2[\sinh px - \sin px], \quad (\text{A.1})$$

where $p = \left(\frac{-\rho\mu_n^2}{\Lambda}\right)^{\frac{1}{4}} > 0$, and $C_i \in \mathbb{R}$. Thereby we already incorporated the zero boundary conditions at $x = 0$. Using the condition $u_n'''(L) = 0$ from (5.14b) yields

$$C_1[\sinh pL - \sin pL] = -C_2[\cosh pL + \cos pL].$$

Since $p \neq 0$ due to $\mu_n \neq 0$, both coefficients are always nonzero. So C_2 can always be determined uniquely from C_1 via this equation. Thus, if (5.14) has a non-trivial solution, it is unique up to multiplicity. This shows that all eigenspaces of \mathcal{B} are one-dimensional, spanned by the Φ_n . Finally, (5.14c) can be used to determine the μ_n for which there is a non-trivial solution. \square

Appendix B. Deferred proofs.

Proof of Lemma 4.3. First we characterize $D(\mathcal{A}^2)$. We use that $y \in D(\mathcal{A}^2)$ if and only if $y \in D(\mathcal{A})$ and $\mathcal{A}y \in D(\mathcal{A})$, or equivalently

$$v \in \tilde{H}_{0,\mathbb{R}}^4(0, L), \quad (\text{B.1})$$

$$u \in \tilde{H}_{0,\mathbb{R}}^6(0, L) \wedge u^{\text{IV}}(0) = u^{\text{V}}(0) = 0, \quad (\text{B.2})$$

$$\xi = Jv'(L), \quad (\text{B.3})$$

$$\psi = mv(L), \quad (\text{B.4})$$

$$u''(L) = \frac{J}{\rho}u^{\text{V}}(L), \quad (\text{B.5})$$

$$\Lambda u'''(L) - k_1(u(L)) - k_2\left(\frac{\psi}{m}\right) = -\frac{m\Lambda}{\rho}u^{\text{IV}}(L). \quad (\text{B.6})$$

It suffices to show that for an arbitrary $y \in D(\mathcal{A})$ we can construct $\{y_n\}_{n \in \mathbb{N}} \subset D(\mathcal{A}^2)$ such that $y_n = [u_n, v_n, \xi_n, \psi_n]^\top$ converges to y in the space $H^4(0, L) \times H^2(0, L) \times \mathbb{R}^2$. Since $\tilde{C}_0^\infty(0, L) := \{f \in C^\infty(0, L) : f^{(k)}(0) = 0, \forall k \in \mathbb{N} \cup \{0\}\}$ is dense in $\tilde{H}_0^2(0, L)$ (see Theorem 3.17 in [1]), there exists a sequence $\{v_n\}_{n \in \mathbb{N}} \subset \tilde{C}_0^\infty(0, L)$ such that $\lim_{n \rightarrow \infty} v_n = v$ in $H^2(0, L)$. Clearly $v_n \in \tilde{H}_{0,\mathbb{R}}^4(0, L)$ for all $n \in \mathbb{N}$. Defining $\xi_n := Jv_n'(L)$ and $\psi_n := mv_n(L)$ ensures that y_n satisfies (B.3) and (B.4). Moreover, the Sobolev embedding $H^2(0, L) \hookrightarrow C^1(0, L)$ implies that $\lim_{n \rightarrow \infty} \xi_n = \xi$ and $\lim_{n \rightarrow \infty} \psi_n = \psi$.

As a final step, we construct a sequence $\{u_n\}_{n \in \mathbb{N}} \subset C^\infty(0, L)$ such that u_n satisfies (B.2), (B.5), and (B.6) for all $n \in \mathbb{N}$, and $\lim_{n \rightarrow \infty} u_n = u$ in $H^4(0, L)$. For this purpose, first we introduce the polynomial $h_n(x) := h_{2,n}x^2 + h_{3,n}x^3 + h_{6,n}x^6 + h_{7,n}x^7 + h_{8,n}x^8 + h_{9,n}x^9 + h_{10,n}x^{10} + h_{11,n}x^{11}$, for $n \in \mathbb{N}$. Next we show that the coefficients $h_{2,n}, \dots, h_{11,n} \in \mathbb{R}$ can uniquely be determined for every $n \in \mathbb{N}$, given certain boundary conditions on h_n . In the following we fix $n \in \mathbb{N}$ arbitrary. From the definition of h_n it is already immediate that

$$h_n(0) = h_n'(0) = h_n^{\text{IV}}(0) = h_n^{\text{V}}(0) = 0, \quad (\text{B.7})$$

holds. Then we set $h_{2,n} = \frac{u''(0)}{2}$ and $h_{3,n} = \frac{u'''(0)}{6}$, which is equivalent to

$$h_n''(0) = u''(0), \quad h_n'''(0) = u'''(0). \quad (\text{B.8})$$

Assume further that

$$h_n^{(k)}(L) = u^{(k)}(L), \quad k \in \{0, 1, 2, 3\},$$

which reads equivalently²:

$$r_1 = h_{n,6} + h_{n,7}L + h_{n,8}L^2 + h_{n,9}L^3 + h_{n,10}L^4 + h_{n,11}L^5 \quad (\text{B.9a})$$

$$r_2 = 6h_{n,6} + 7h_{n,7}L + 8h_{n,8}L^2 + 9h_{n,9}L^3 + 10h_{n,10}L^4 + 11h_{n,11}L^5 \quad (\text{B.9b})$$

$$\begin{aligned} r_3 = 6^2h_{n,6} + 7^2h_{n,7}L + 8^2h_{n,8}L^2 + 9^2h_{n,9}L^3 \\ + 10^2h_{n,10}L^4 + 11^2h_{n,11}L^5 \end{aligned} \quad (\text{B.9c})$$

$$\begin{aligned} r_4 = 6^3h_{n,6} + 7^3h_{n,7}L + 8^3h_{n,8}L^2 + 9^3h_{n,9}L^3 \\ + 10^3h_{n,10}L^4 + 11^3h_{n,11}L^5, \end{aligned} \quad (\text{B.9d})$$

²Here we use the notation $k^l := k!/(k-l)!$ for $k, l \in \mathbb{N}$, $k \geq l$.

where

$$\begin{aligned} r_1 &= \frac{u(L)}{L^6} - \frac{u''(0)}{2L^4} - \frac{u'''(0)}{6L^3}, & r_2 &= \frac{u'(L)}{L^5} - \frac{u''(0)}{L^4} - \frac{u'''(0)}{2L^3}, \\ r_3 &= \frac{u''(L)}{L^4} - \frac{u''(0)}{L^4} - \frac{u'''(0)}{L^3}, & r_4 &= \frac{u'''(L)}{L^3} - \frac{u'''(0)}{L^3}. \end{aligned}$$

Finally the two additional conditions are imposed on h_n :

$$\frac{m\Lambda}{\rho} h_n^{\text{IV}}(L) = -\Lambda u'''(L) + k_1(u(L)) + k_2\left(\frac{\psi_n}{m}\right), \quad (\text{B.10})$$

$$\frac{J}{\rho} h_n^{\text{V}}(L) = u''(L). \quad (\text{B.11})$$

(B.10) and (B.11) are equivalent to:

$$6^4 h_{n,6} + 7^4 h_{n,7} L + 8^4 h_{n,8} L^2 + 9^4 h_{n,9} L^3 + 10^4 h_{n,10} L^4 + 11^4 h_{n,11} L^5 = r_5, \quad (\text{B.12a})$$

$$6^5 h_{n,6} + 7^5 h_{n,7} L + 8^5 h_{n,8} L^2 + 9^5 h_{n,9} L^3 + 10^5 h_{n,10} L^4 + 11^5 h_{n,11} L^5 = r_6, \quad (\text{B.12b})$$

with

$$r_5 = \rho \frac{-\Lambda u'''(L) + k_1(u(L)) + k_2\left(\frac{\psi_n}{m}\right)}{\Lambda m L^2}, \quad r_6 = \frac{\rho u''(L)}{JL}.$$

The linear system consisting of (B.9) and (B.12) has a strictly positive determinant. Hence, its solution h_n exists and is unique. Consequently, (B.7), (B.8), and (B.9) imply that $u - h_n \in H_0^4(0, L)$, for all $n \in \mathbb{N}$. Since $C_0^\infty(0, L)$ is dense in $H_0^4(0, L)$, there exists a sequence $\{\tilde{u}_n\}_{n \in \mathbb{N}} \subset C_0^\infty(0, L)$ such that $\|\tilde{u}_n - (u - h_n)\|_{H^4} < \frac{1}{n}$, $\forall n \in \mathbb{N}$. Now defining $u_n = \tilde{u}_n + h_n$, gives $\lim_{n \rightarrow \infty} u_n = u$ in $H^4(0, L)$. Obviously u_n satisfies (B.2) for all $n \in \mathbb{N}$. Also due to (B.10) and (B.11), u_n satisfies (B.5) and (B.6), as well. The statement follows. \square

Proof of Lemma 5.4. We first consider $y_0 \in D(\mathcal{A})$. Then $S(t)y_0$ is the classical solution of (3.2), and satisfies the integrated equation:

$$S(t)y_0 - y_0 = \int_0^t \mathcal{A}S(s)y_0 \, ds + \int_0^t \mathcal{N}S(s)y_0 \, ds.$$

Since $S(t)y_0 \in C^1(\mathbb{R}^+; \mathcal{H})$ and \mathcal{N} is locally Lipschitz continuous, we find that both $t \mapsto \mathcal{N}S(t)y_0$ and $t \mapsto \mathcal{A}S(t)y_0$ are continuous, so $\mathcal{A}S(t)y_0 \in C(\mathbb{R}^+; \mathcal{H})$. Therefore we may write for any $t > 0$:

$$\begin{aligned} \int_0^t S(s)y_0 \, ds &= \lim_{N \rightarrow \infty} \sum_{j=1}^N \frac{t}{N} S\left(\frac{jt}{N}\right)y_0, \\ \int_0^t \mathcal{A}S(s)y_0 \, ds &= \lim_{N \rightarrow \infty} \sum_{j=1}^N \frac{t}{N} \mathcal{A}S\left(\frac{jt}{N}\right)y_0 = \lim_{N \rightarrow \infty} \mathcal{A} \sum_{j=1}^N \frac{t}{N} S\left(\frac{jt}{N}\right)y_0, \end{aligned}$$

due to the linearity of \mathcal{A} . Since \mathcal{A} is skew-adjoint, see Proposition A.1, it is closed. So we obtain that

$$\int_0^t S(s)y_0 \, ds \in D(\mathcal{A}), \quad \mathcal{A} \int_0^t S(s)y_0 \, ds = \int_0^t \mathcal{A}S(s)y_0 \, ds.$$

So there holds (5.3) for $y_0 \in D(\mathcal{A})$ and any $t > 0$.

Let now $y_0 \in \mathcal{H} \setminus D(\mathcal{A})$, and $\{y_{n,0}\} \subset D(\mathcal{A})$ such that $y_{n,0} \rightarrow y_0$ as $n \rightarrow \infty$. For every $T > 0$ we have $S(t)y_{n,0} \rightarrow S(t)y_0 \in C([0, T], \mathcal{H})$. Since furthermore \mathcal{N} is locally Lipschitz continuous, we get for every $t > 0$:

$$\begin{aligned} \lim_{n \rightarrow \infty} (S(t)y_{n,0} - y_{n,0}) &= (S(t)y_0 - y_0), \\ \lim_{n \rightarrow \infty} \int_0^t \mathcal{N}S(s)y_{n,0} \, ds &= \int_0^t \mathcal{N}S(s)y_0 \, ds. \end{aligned}$$

Applying those two limits in (5.3) for $y_{n,0}$ we obtain:

$$\lim_{n \rightarrow \infty} A \int_0^t S(s)y_{n,0} \, ds = S(t)y_0 - y_0 - \int_0^t \mathcal{N}S(s)y_0 \, ds.$$

But there also holds

$$\lim_{n \rightarrow \infty} \int_0^t S(s)y_{n,0} \, ds = \int_0^t S(s)y_0 \, ds.$$

Since A is closed, these last two limits prove (5.2) and (5.3) for $y_0 \in \mathcal{H} \setminus D(\mathcal{A})$. \square

Proof of Lemma 5.6. The general solution $\varphi \in \tilde{H}_0^4(0, L)$ to (5.14a) is of the form (A.1), with $p = \left(\frac{-\rho\mu^2}{\Lambda}\right)^{\frac{1}{4}} > 0$. The boundary conditions (5.14b) and (5.14c) are now equivalent to the following two equations for C_1 and C_2 :

$$C_1 (\sinh pL - \sin pL) + C_2 (\cosh pL + \cos pL) = 0, \quad (\text{B.13})$$

and

$$\begin{aligned} C_1 [J\mu^2 (\sinh pL + \sin pL) + p\Lambda (\cosh pL + \cos pL)] \\ + C_2 [J\mu^2 (\cosh pL - \cos pL) + p\Lambda (\sinh pL + \sin pL)] = 0. \end{aligned} \quad (\text{B.14})$$

Furthermore, the additional condition $\varphi(L) = 0$ reads

$$C_1 (\cosh pL - \cos pL) + C_2 (\sinh pL - \sin pL) = 0. \quad (\text{B.15})$$

First we use (B.13) and (B.15) to determine the constants C_1 and C_2 . In order for φ to be non-zero the determinant of the linear system formed by (B.13) and (B.15) needs to vanish, i.e.

$$\begin{aligned} (\sinh pL - \sin pL)^2 - (\cosh pL - \cos pL)(\cosh pL + \cos pL) \\ = -2 \sinh pL \sin pL = 0. \end{aligned}$$

Since $pL > 0$, this is true iff $p = \frac{\ell\pi}{L}$ for some $\ell \in \mathbb{N}$. Hence $\mu^2 = -\frac{\Lambda}{\rho} \left(\frac{\ell\pi}{L}\right)^4$. Now (B.13) gives $C_2 = -C_1 \frac{\sinh \ell\pi}{\cosh \ell\pi + (-1)^\ell}$. Now we investigate in which situation also the third condition (B.14) is fulfilled. Multiplying (B.14) by $\frac{(-1)^\ell \cosh \ell\pi + 1}{2C_1}$, we get

$$-J \frac{\Lambda}{\rho} \left(\frac{\ell\pi}{L}\right)^4 \sinh \ell\pi + \frac{\ell\pi\Lambda}{L} [\cosh \ell\pi + (-1)^\ell] = 0,$$

and equivalently

$$J = \rho \left(\frac{L}{\ell\pi}\right)^3 \frac{\cosh \ell\pi + (-1)^\ell}{\sinh \ell\pi}.$$

In this case, the eigenfunction φ is given by (up to normalization)

$$\varphi(x) = \left(\cosh \frac{\ell\pi x}{L} - \cos \frac{\ell\pi x}{L} \right) - \frac{\sinh \ell\pi}{\cosh \ell\pi + (-1)^\ell} \left(\sinh \frac{\ell\pi x}{L} - \sin \frac{\ell\pi x}{L} \right). \quad (\text{B.16})$$

\square

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