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# Control of a Wave Equation with a Dynamic Boundary Condition 

Nicolas Vanspranghe, Francesco Ferrante, Christophe Prieur


#### Abstract

The general problem of this paper is the analysis of wave propagation in a bounded medium where the uncontrolled boundary obeys a coupled differential equation. More precisely, we study a one-dimensional wave equation with a nonlinear second-order dynamic boundary condition and a Neuman-type boundary control acting on the other extremity. A generic class of nonlinear collocated feedback laws is considered. Hadamard well-posedness is established for the closed-loop system, with initial data lying in the natural energy space of the problem. Moreover, we investigate an example of stabilization through a proportional controller.


## I. Introduction

The aim of this paper is to study a wave equation in a bounded one-dimensional medium supplied with a dynamic (or kinetic) boundary condition at one end. The system is actuated via a Neuman-type control at the other end. Dynamic boundary conditions involve second-order time derivative and are typically obtained in physical models for which the momentum of the boundary cannot be neglected. A prime example is an infinite-dimensional model for the propagation of mechanical vibrations along drilling rods. In that case, the control is the torque applied to one extremity, and the kinetic boundary condition is given by the behavior of the rock-tip contact, which is subject to nonlinear friction. In particular, stick-slip phenomena can occur at the rock-tip interface and generate unwanted vibrations that might jeopardize the plant. A mechanical analysis of the rock-tip dynamics is given in [7]. The problem of stabilization and regulation of the velocity at the rock-tip contact has sparked interest in the control commmunity, see e.g. [10]. In engineering applications, the goal is to minimize the stick-slip vibrations through a suitable control. Various boundary control strategies are proposed to address this problem in [8], including backstepping design - see also [9]. In [4], an observer-based boundary control design is proposed. In [13], stabilization and regulation using a proportional integral boundary controller is investigated. In [12] and [1], a similar problem is considered, but with a boundary anti-damping only involving first-order derivatives. However, the aforementioned papers deal with linearized models. We should also mention [3] and [6] where different classes of boundary nonlinearities for distributed parameter systems are considered.

Let us introduce the specific dynamical model under study in this paper. Let $L>0$ and $\Omega \triangleq(0, L) \subset \mathbb{R}$. We deal with

[^0]the following control system:
\[

$$
\begin{cases}\partial_{t t} u-\partial_{x x} u=0 & \text { on } \Omega \times \mathbb{R}^{+}  \tag{1}\\ \partial_{t t} u(0, t)-\partial_{x} u(0, t)=F\left(\partial_{t} u(0, t)\right) & \text { for all } t \\ \partial_{x} u(L, t)=U(t) & \text { for all } t\end{cases}
$$
\]

where $F$ is a scalar function that models the nonlinear boundary friction, and $U(t)$ is the control input.

Our first goal is to prove the well-posedness of the control system (1) supplied with the following collocated feedback:

$$
\begin{equation*}
U(t)=-g\left(\partial_{t} u(L, t)\right) \tag{2}
\end{equation*}
$$

where $g$ is a continuous nondecreasing scalar function. The output considered here is the velocity at the boundary where the actuator lies, meaning that an observer is not required to implement the controller. Note that this class of feedback laws includes nonlinearities of interest in control applications, such as saturations or deadzones - see e.g. [5].

In this paper, we prove that, under appropriate assumptions, closing the loop in (1) with (2) leads to a wellposed dynamical system. A precise variational formulation of the problem is given and the regularity of the solutions is rigorously investigated. The underlying control problem is the stabilization of a possibly non-dissipative boundary by an action on the other boundary. We prove that, under suitable conditions, exponential stability can be achieved using a proportional controler. To the best of our knowledge, the stabilization of such system in the presence of nonlinear boundary anti-damping has not been investigated so far.

This paper is organized as follows. Section $I I$ introduces the variational formulation of control system (1) and contains the first main result, namely the well-posedness of the closedloop dynamics. Section III introduces the natural energy associated with (1) and states an exponential stability result under appropriate assumptions on the nonlinear map $F$. This is the second main contribution. Section IV contains some illustrative numerical results. Section V gives the proof of the well-posedness result. Section VI contains concluding remarks.

Notation: Given a Banach space $E$, we denote its norm by $\|\cdot\|_{E}$ and we use the duality bracket $\langle\phi, x\rangle_{E}$ to write $\phi(x)$ for any $x \in E$ and $\phi \in E^{\prime}$, where $E^{\prime}$ is the topological dual of $E$. If $E$ is also a Hilbert space, its scalar product is written $(\cdot, \cdot)_{E}$. The space of infinitely differentiable functions on $\Omega$ with compact support is denoted by $\mathcal{C}_{c}^{\infty}(\Omega)$. Also, for $T>0$, we denote by $W^{1, p}(0, T ; E)$ the subspace of $L^{p}(0, T ; E)$ composed of (classes of) $E$-valued functions $f$ such that, for some $g$ in $L^{p}(0, T ; E), f(t)=f(0)+\int_{0}^{t} g(s) \mathrm{d} s$ for a.e. $t$ in $(0, T)$. Such class $f$ is identified with its continuous
representative and we say that $f^{\prime}=g$ in the sense of $E$ valued distributions.

## II. VARIATIONAL FORMULATION AND WELL-POSEDNESS

In this section, we establish the framework in which we analyze system (1), and state our well-posedness result. We start by introducing the energy spaces associated with (1) as well as some notation. First, define

$$
\begin{equation*}
H \triangleq L^{2}(\Omega) \times \mathbb{R} \tag{3}
\end{equation*}
$$

We endow $H$ with the usual product Hilbertian structure. Define now the following subspace of $H$ :

$$
\begin{equation*}
V \triangleq\left\{(u, u(0)): u \in H^{1}(\Omega)\right\} \tag{4}
\end{equation*}
$$

which is equipped with the scalar product inherited from $H^{1}(\Omega) \times \mathbb{R}$. As stated in Section V, $V$ is also a Hilbert space. We also introduce the state space

$$
\begin{equation*}
\mathcal{X} \triangleq V \times H \tag{5}
\end{equation*}
$$

endowed with the product Hilbertian structure. For the sake of clarity, we use parenthesis to denote elements of $V$ or $H$, and brackets to denote elements of $\mathcal{X}$, as in $\mathbf{X}=[\mathbf{u}, \mathbf{v}] \in \mathcal{X}$, $\mathbf{u}=(u, u(0)) \in V$, etc. Now, let us define the bilinear continuous symmetric form $a$ on $V \times V$ by

$$
\begin{equation*}
a\left(\mathbf{u}_{1}, \mathbf{u}_{2}\right) \triangleq \int_{\Omega} \partial_{x} u_{1}(x) \partial_{x} u_{2}(x) \mathrm{d} x \tag{6}
\end{equation*}
$$

Finally, we denote by $\delta_{L}$ the linear form mapping $\mathbf{u} \in V$ into $u(L)$, which belongs to $V^{\prime}$ since $H^{1}(\Omega)$ is continuously embedded into $\mathcal{C}(\bar{\Omega})$ in the one-dimensional case.

Assumption 1. The scalar function $F$ is globally Lipschitz.
We define the nonlinear operator $B$ on $H$ associated with the first-order boundary pertubation:

$$
\begin{equation*}
\forall \mathbf{v}=(v, \gamma) \in H, \quad B(\mathbf{v}) \triangleq(0, F(\gamma)) \tag{7}
\end{equation*}
$$

This operator is globally Lipschitz by Assumption 1. Let us precise the meaning of weak solutions to (1).

Definition 1. A weak solution to (1) on $(0, T)$ is any $\mathbf{u}$ in $L^{\infty}(0, T ; V) \cap W^{1, \infty}(0, T ; H)$ verifying

$$
\begin{align*}
\int_{0}^{T} & -\left(\mathbf{u}^{\prime}(t), \phi^{\prime}(t)\right)_{H}+a(\mathbf{u}(t), \phi(t)) \mathrm{d} t=  \tag{8}\\
& \int_{0}^{T}\left(B\left(\mathbf{u}^{\prime}(t)\right), \phi(t)\right)_{H}+U(t)\left\langle\delta_{L}, \boldsymbol{\phi}(t)\right\rangle_{V} \mathrm{~d} t
\end{align*}
$$

for all test-functions $\phi$ in $\mathcal{C}_{c}(0, T ; V) \cap \mathcal{C}_{c}^{1}(0, T ; H)$.
Definition 1 is motivated by simple formal calculations in which one multiplies the wave equation by some smooth test-function $\phi(x, t)$ and integrates it over $\Omega \times(0, T)$ using integration by parts and the boundary conditions. Note that Definition 1 makes sense if, say, $U$ belongs to $L^{2}(0, T)$. Closing the loop, a hidden regularity property of the solutions is needed to ensure all terms are defined ${ }^{1}$

[^1]Assumption 2. The function $g: \mathbb{R} \rightarrow \mathbb{R}$ is continuous and nondecreasing. ${ }^{2}$

We show that the closed-loop system generates a dynamical system on $\mathcal{X}$ by defining the operators $\mathcal{S}_{t}$ for $t \geq 0$ as follows:

$$
\begin{equation*}
\forall \mathbf{X}=\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right] \in \mathcal{X}, \quad \mathcal{S}_{t}(\mathbf{X}) \triangleq\left[\mathbf{u}(t), \mathbf{u}^{\prime}(t)\right] \in \mathcal{X} \tag{9}
\end{equation*}
$$

where $\mathbf{u}$ is the unique solution associated with initial data $\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right]$.
Theorem 1 (Hadamard well-posedness). Let $\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right] \in \mathcal{X}$. Under Assumptions 1 and 2 there exists a unique (weak) solution $\mathbf{u} \in \mathcal{C}\left(\mathbb{R}^{+}, V\right) \cap \mathcal{C}^{1}\left(\mathbb{R}^{+}, H\right)$ to (1) with feedback (2) and initial data $\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right]$. The solution $\mathbf{u}$ enjoys the following hidden regularity property: for all $T>0$,

$$
\begin{equation*}
u(L, \cdot) \in H^{1}(0, T) \text { and } g\left(\partial_{t} u(L, \cdot)\right) \in L^{2}(0, T) \tag{10}
\end{equation*}
$$

Moreover, we can associate with (1) and control law (2) the semigroup $\mathfrak{S}=\left\{\mathcal{S}_{t}\right\}_{t \geq 0}$ of nonlinear continuous operators on $\mathcal{X}$ as defined in (9).

Theorem 1 is proved in Section V. In the following proposition, we give some additional regularity when the initial datum is smooth and verifies a compatibility condition. We write $W \triangleq\left[H^{2}(\Omega) \times \mathbb{R}\right] \cap V$, equipped with the scalar product inherited from $H^{2}(\Omega) \times \mathbb{R}$.
Proposition 1 (Strong solutions). Let $T>0$. If $\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right]$ belongs to $W \times V$ and verifies $\partial_{x} u^{0}(L)=-g\left(u^{1}(L)\right)$, then,

$$
\begin{equation*}
\mathbf{u} \in L^{\infty}(0, T ; W), \quad \mathbf{u}^{\prime} \in L^{\infty}(0, T ; V) \tag{11}
\end{equation*}
$$

Moreover, any weak solution $\mathbf{u}$ is the limit of a sequence of strong solutions $\left[\mathbf{u}_{n}, \mathbf{u}_{n}^{\prime}\right]$ in $\mathcal{C}([0, T], \mathcal{X})$, and $\partial_{t} u_{n}(L, \cdot) \rightarrow$ $\partial_{t} u(L, \cdot)$ in $L^{2}(0, T)$.

Proposition 1 is a byproduct of the proof of Theorem 1 and is used to justify computations performed in Section III

## III. Stability analysis of the Closed-LOOP System

Next, we analyze the stability of the closed-loop system when the feedback is linear. We introduce the energy functional $\mathcal{E}$ defined on $\mathcal{X}$ by

$$
\begin{equation*}
\mathcal{E}(\mathbf{u}, \mathbf{v}) \triangleq \frac{1}{2}\left[\|\mathbf{v}\|_{H}^{2}+a(\mathbf{u}, \mathbf{u})\right] \tag{12}
\end{equation*}
$$

which is, in the context of abstract wave equations, the natural "mechanical" energy. Here, for all $\mathbf{u}$ in $V$ and $\mathbf{v}=(v, \gamma)$ in $H$,

$$
\begin{equation*}
\mathcal{E}(\mathbf{u}, \mathbf{v})=\frac{1}{2} \int_{\Omega}|v(x)|^{2}+\left|\partial_{x} u(x)\right|^{2} \mathrm{~d} x+\frac{1}{2}|\gamma|^{2} \tag{13}
\end{equation*}
$$

For any $\rho$ and $\mu$ in $H^{1}(\Omega)$, we define a weighted energy functional (or Lyapunov function candidate):

$$
\Gamma_{\rho, \mu}(\mathbf{u}, \mathbf{v}) \triangleq \frac{1}{2} \int_{\Omega}\left[\begin{array}{c}
v  \tag{14}\\
\partial_{x} u
\end{array}\right]^{\top}\left[\begin{array}{cc}
\mu & \rho \\
\rho & \mu
\end{array}\right]\left[\begin{array}{c}
v \\
\partial_{x} u
\end{array}\right]+\frac{1}{2} \mu(0)|\gamma|^{2}
$$

[^2]It follows from the boundedness of the weights $\rho$ and $\mu$ that there exists a constant $M_{\rho, \mu}>0$ such that

$$
\begin{equation*}
\forall[\mathbf{u}, \mathbf{v}] \in \mathcal{X}, \quad \Gamma_{\rho, \mu}(\mathbf{u}, \mathbf{v}) \leq M_{\rho, \mu} \mathcal{E}(\mathbf{u}, \mathbf{v}) \tag{15}
\end{equation*}
$$

A sufficient condition for the converse inequality to hold is the existence of $C_{1}, C_{2}>0$ such that

$$
\begin{equation*}
\mu(x) \geq C_{1} \text { and } \mu(x)^{2}-\rho(x)^{2} \geq C_{2} \tag{16}
\end{equation*}
$$

holds for all $x$ in $\Omega$. We also remark that $\mathcal{E}$ and $\Gamma$ are continuous on $\mathcal{X}$. We note that the position does not appear in the energy functionals and without Poincaré-type inequalities, one cannot directly infer an estimate of the norm of the position. Stabilization shall be investigated with respect to the following invariant set, which is the line composed of constant solutions:

$$
\begin{equation*}
\mathfrak{A} \triangleq\{[\theta \mathbb{1}, 0] \in \mathcal{X}: \theta \in \mathbb{R}\} \tag{17}
\end{equation*}
$$

where $\mathbb{1} \triangleq((x \mapsto 1), 1)$. Observe that $\mathfrak{A}$ is exactly the subset of $\mathcal{X}$ where the mechanical energy $\mathcal{E}$ vanishes. For any $[\mathbf{u}, \mathbf{v}]$ in $\mathcal{X}$,

$$
\begin{equation*}
\operatorname{dist}([\mathbf{u}, \mathbf{v}], \mathfrak{A})^{2} \leq C \mathcal{E}(\mathbf{u}, \mathbf{v}) \tag{18}
\end{equation*}
$$

where $C$ is some positive constant. This is a consequence of the Poincaré-Wirtinger inequality - see [2] for instance. From now on, we are interested in stabilizing system (1) using a proportional feedback, i.e.

$$
\begin{equation*}
U(t)=-\alpha \partial_{t} u(L, t) \tag{19}
\end{equation*}
$$

where $\alpha$ is a positive gain to be tuned.
Assumption 3. The function $F$ is $q$-Lipschitz for some $q<$ $1 / 2$. Also, $F(0)=0$.

The condition $F(0)=0$ is quite natural when dealing with friction. Assumption 3 is meant to be in force in the context of nonlinear anti-damping ${ }^{3}$

Theorem 2 (Exponential stability). Under Assumption 3 and with $\alpha=1, \mathfrak{A}$ uniformly and exponentially attracts the bounded sets of $\mathcal{X}$. More precisely, there exist two positive constants $M$ and $\eta$ such that for all solutions $\mathbf{u}$ to (1),

$$
\begin{equation*}
\mathcal{E}\left(\mathbf{u}(t), \mathbf{u}^{\prime}(t)\right) \leq M \mathcal{E}\left(\mathbf{u}^{0}, \mathbf{u}^{1}\right) \exp (-\eta t) \tag{20}
\end{equation*}
$$

Moreover, for each solution $\mathbf{u}, \mathbf{u}(t)$ converges in $V$ to some constant function $\mathbf{u}_{\infty}$ when $t \rightarrow+\infty$.

Proof. We pick a solution $\mathbf{u}$ and, for the sake of concision, we denote by $\Gamma_{\rho, \mu}^{\mathbf{u}}$ the function $t \in \mathbb{R}^{+} \mapsto \Gamma_{\rho, \mu}\left(\mathbf{u}(t), \mathbf{u}^{\prime}(t)\right)$, which is continuous. Take $\rho, \mu \in H^{1}(\Omega)$. Assume temporarily that $\mathbf{u}$ is a strong solution. For all $\tau \geq 0$, denoting by

[^3]$Q_{\tau}$ the rectangle $\Omega \times(0, \tau)$, we have the following identity:
\[

$$
\begin{align*}
&\left.\Gamma_{\rho, \mu}^{\mathbf{u}}(t)\right|_{0} ^{\tau}=-\frac{1}{2} \iint_{Q_{\tau}}\left[\begin{array}{c}
\partial_{t} u \\
\partial_{x} u
\end{array}\right]^{\top}\left[\begin{array}{cc}
\rho^{\prime} & \mu^{\prime} \\
\mu^{\prime} & \rho^{\prime}
\end{array}\right]\left[\begin{array}{c}
\partial_{t} u \\
\partial_{x} u
\end{array}\right] \\
&+\frac{1}{2}\left[\rho(L)\left(1+\alpha^{2}\right)-2 \alpha \mu(L)\right] \int_{0}^{\tau}\left|\partial_{t} u(L, t)\right|^{2} \mathrm{~d} t  \tag{21}\\
&-\frac{\rho(0)}{2} \int_{0}^{\tau}\left|\partial_{t} u(0, t)\right|^{2}+\left|\partial_{x} u(0, t)\right|^{2} \mathrm{~d} t \\
&+\mu(0) \int_{0}^{\tau} F\left(\partial_{t} u(0, t)\right) \partial_{t} u(0, t) \mathrm{d} t
\end{align*}
$$
\]

Equation 21) is obtained by multiplying the wave equation $\partial_{t t} u(x, t)-\partial_{x x} u(x, t)=0$ by $\phi(x, y)=\mu(x) \partial_{t} u(x, t)+$ $\rho(x) \partial_{x} u(x, t) \in H^{1}\left(Q_{\tau}\right)$, integrating over $Q_{\tau}$, and performing a few integrations by parts. If $\rho$ and $\mu$ are nonnegative, writing $|F(s)| \leq q|s|$ where $q$ is the Lipschitz constant in Assumption 3, we have the following inequality:

$$
\begin{align*}
&\left.\Gamma_{\rho, \mu}^{\mathbf{u}}(t)\right|_{0} ^{\tau} \leq-\frac{1}{2} \iint_{Q_{\tau}}\left[\begin{array}{c}
\partial_{t} u \\
\partial_{x} u
\end{array}\right]^{\top}\left[\begin{array}{cc}
\rho^{\prime} & \mu^{\prime} \\
\mu^{\prime} & \rho^{\prime}
\end{array}\right]\left[\begin{array}{c}
\partial_{t} u \\
\partial_{x} u
\end{array}\right] \\
&+\frac{1}{2}\left[\rho(L)\left(1+\alpha^{2}\right)-2 \alpha \mu(L)\right] \int_{0}^{\tau}\left|\partial_{t} u(L, t)\right|^{2} \mathrm{~d} t  \tag{22}\\
&+\frac{1}{2}[2 q \mu(0)-\rho(0)] \int_{0}^{\tau}\left|\partial_{t} u(0, t)\right|^{2} \mathrm{~d} t
\end{align*}
$$

Take $\mu(x)=1$. For (16) to hold, it suffices to have $\rho(x) \leq$ $1-\epsilon$ for some $\epsilon>0$. Now, we derive some sufficient conditions for the energy to decay exponentially. It suffices to have $\rho(0) \geq 2 q+\epsilon$, and $\rho^{\prime}(x) \geq \epsilon$ a.e. as well as

$$
\begin{equation*}
\rho(L) \leq \frac{2 \alpha}{1+\alpha^{2}}-\epsilon \tag{23}
\end{equation*}
$$

for some $\epsilon>0$. Since $q<1 / 2$, there exists an increasing affine function $\rho$ such that $\rho(0)>2 q$ and $\rho(L)<1$. Let $\alpha=1$ so that 23 holds. As a result, by Grönwall's lemma, we obtain the following: with this particular choice of $\rho$ and $\mu$, there exists a positive constant (solution independent) $\eta$ such that

$$
\begin{equation*}
\forall t \geq 0, \quad \Gamma_{\rho, \mu}^{\mathbf{u}}(t) \leq \Gamma_{\rho, \mu}^{\mathbf{u}}(0) \exp (-\eta t) \tag{24}
\end{equation*}
$$

By a density-continuity argument, the uniform estimate 24, holds for weak solutions as well - see Proposition 1 Since (16) holds, then there exists $M>0$ such that (20) is verified by any solution. As for the second statement of Theorem 2, let $\mathbf{u}$ be a solution. Take an increasing sequence of nonnegative real numbers $t_{n}$ such that $t_{n} \rightarrow+\infty$ when $n \rightarrow+\infty$. Then $\left\{\mathbf{u}\left(t_{n}\right)\right\}_{n \geq 0}$ is a Cauchy sequence in $V$. Indeed, for any $m \geq n$, writing the variation of $\mathbf{u}(t)$ between $t_{n}$ and $t_{m}$, we have

$$
\begin{align*}
\left\|\mathbf{u}\left(t_{m}\right)-\mathbf{u}\left(t_{n}\right)\right\|_{H} & \leq \int_{t_{n}}^{t_{m}}\left\|\mathbf{u}^{\prime}(s)\right\|_{H} \mathrm{~d} s  \tag{25}\\
& \leq M^{\prime} \int_{t_{n}}^{+\infty} \exp (-\eta s / 2) \mathrm{ds}
\end{align*}
$$

which converges to 0 when $n \rightarrow+\infty$; also, we already know that $a(\mathbf{u}(t), \mathbf{u}(t)) \rightarrow 0$ when $t \rightarrow+\infty$, thus $\| \mathbf{u}\left(t_{m}\right)-$ $\mathbf{u}\left(t_{n}\right) \|_{V}$ can be arbitrarily small. Using a similar argument,
one verifies that the limit does not depend on the sequence $\left\{t_{n}\right\}$. Therefore, $\mathbf{u}(t)$ converges to some $\mathbf{u}_{\infty} \in V$ and $\mathbf{u}_{\infty}$ is constant by (18).

## IV. Numerical simulations

We provide some numerical computations for illustrative purposes.


Fig. 1. Evolution of the boundary position $u(0, t)$ over time. It obeys a second-order differential equation coupled with the wave equation.


Fig. 2. Evolution of the mechanical energy $\mathcal{E}\left(\mathbf{u}(t), \mathbf{u}^{\prime}(t)\right)$ over time. The uncontrolled system would be unstable due to the boundary anti-damping.

We discretize (1) using finite elements over space and finite differences over time, on the basis of the functional formulation of the problem. Figures 1 and 2 are obtained with the following parameters: we take $\alpha=1, F(x)=q x$ with $q=0.1$ and $L=1$. The domain $\Omega$ is discretized into 100 points and the time step is set to 0.001 . Further computations suggest that the condition on $q$ derived in Theorem 2 is nearly sharp as taking $q=0.5$ leads to an exponentially unstable system. Computations also suggest that the proportional feedback is not robust to (numerical) errors: taking values of $q$ slightly below 0.5 induces unclear situations where an exponential decay is not easily identifiable.

## V. Proof of the Hadamard theorem

The proof of Theorem 1 relies on nonlinear semigroup techniques and appropriate energy estimates. We begin this
section with some remarks on the functional settings of the problem. The following result states some useful properties of the spaces $H$ and $V$.

Lemma 1. $V$ is a separable Hilbert space isomorphic to $H^{1}(\Omega)$. Moreover, $V$ is a dense subset of $H$.

Since $\|\mathbf{u}\|_{V}^{2}=\|\mathbf{u}\|_{H}^{2}+a(\mathbf{u}, \mathbf{u})$, we see that $V$ is continuously embedded into $H$. We denote $H^{\prime}$ the topological dual of $H$. By Riesz representation theorem, we make the identification $H \simeq H^{\prime}$. Because $V$ is a dense subset of $H$, the latter can be identified as a dense subset of $V^{\prime}$, leading to the classical injection chain $V \hookrightarrow H \simeq H^{\prime} \hookrightarrow V^{\prime}$, where each space is dense and continuously embedded into the following one. We denote by $A$ the continuous operator from $V$ into $V^{\prime}$ defined by $\left\langle A \mathbf{u}_{1}, \mathbf{u}_{2}\right\rangle_{V} \triangleq a\left(\mathbf{u}_{1}, \mathbf{u}_{2}\right)$ for all $\mathbf{u}_{1}, \mathbf{u}_{2}$ in $V$. Next, define an unbounded (nonlinear) operator $\mathcal{A}_{g}$ on $\mathcal{X}$ by

$$
\begin{align*}
& \mathcal{D}\left(\mathcal{A}_{g}\right) \triangleq\left\{[\mathbf{u}, \mathbf{v}] \in W \times V: \partial_{x} u(L)=-g(v(L))\right\} \\
& \forall \mathbf{X}=[\mathbf{u}, \mathbf{v}] \in \mathcal{D}\left(\mathcal{A}_{g}\right), \quad \mathcal{A}_{g}(\mathbf{X}) \triangleq-\left[\begin{array}{c}
\mathbf{v} \\
\left(\partial_{x x} u, \partial_{x} u(0)\right)
\end{array}\right] \tag{26}
\end{align*}
$$

Note that the domain $\mathcal{D}\left(\mathcal{A}_{g}\right)$ need not be a subspace. We start with the following first-order abstract Cauchy problem:

$$
\left\{\begin{array}{l}
\dot{\mathbf{X}}(t)+\mathcal{A}_{g}(\mathbf{X}(t))=\mathcal{F}(\mathbf{X}(t))  \tag{27}\\
\mathbf{X}(0)=\mathbf{X}^{0}
\end{array}\right.
$$

where $\mathcal{F}$ is the nonlinear perturbation operator on $\mathcal{X}$ defined by $\mathcal{F}(\mathbf{X}) \triangleq[0, B(\mathbf{v})]$ for all $\mathbf{X}=[\mathbf{u}, \mathbf{v}] \in \mathcal{X}$. We see that $\mathcal{F}$ is Lipschitz, since the $H$-valued mapping $B$ is Lipschitz. We wish to prove that (27) is a Lipschitz perturbation of an evolution equation with maximal monotone generator, hence the following result.

Proposition 2. The unbounded operator $\mathcal{A}_{g}+\mathrm{id}$ is maximal monotone.
Proof. We start with the monotonicity, and then we shall prove the surjectivity of $\left[\mathcal{A}_{g}+\mathrm{id}\right]+\mathrm{id}$.
Step 1: Monotonicity. Let $\mathbf{X}_{1}=\left[\mathbf{u}_{1}, \mathbf{v}_{1}\right]$ and $\mathbf{X}_{2}=$ $\left[\mathbf{u}_{2}, \mathbf{u}_{2}\right]$ in $\mathcal{D}\left(\mathcal{A}_{g}\right)$. We denote $\mathbf{u}_{1}-\mathbf{u}_{2}$ (resp. $\left.\mathbf{v}_{1}-\mathbf{v}_{2}\right)$ by $\mathbf{u}$ (resp. v), and also $\mathbf{X}_{1}-\mathbf{X}_{2}$ by $\mathbf{X}$. We have

$$
\begin{align*}
\left(\mathcal{A}_{g}\left(\mathbf{X}_{1}\right)-\right. & \left.\mathcal{A}_{g}\left(\mathbf{X}_{2}\right), \mathbf{X}\right)_{\mathcal{X}}=-(\mathbf{u}, \mathbf{v})_{V}  \tag{28}\\
& -\left(\partial_{x x} u, v\right)_{L^{2}(\Omega)}-\partial_{x} u(0) v(0)
\end{align*}
$$

Recall that $(\mathbf{u}, \mathbf{v})_{V}=(\mathbf{u}, \mathbf{v})_{H}+a(\mathbf{u}, \mathbf{v})$. Moreover, by integration by parts, we have $a(\mathbf{u}, \mathbf{v})=-\left(\partial_{x x} u, v\right)_{L^{2}(\Omega)}+$ $\partial_{x} u(L) v(L)-\partial_{x} u(0) v(0)$. Thus, from 28 we obtain

$$
\begin{equation*}
\left(\mathcal{A}_{g}\left(\mathbf{X}_{1}\right)-\mathcal{A}_{g}\left(\mathbf{X}_{2}\right), \mathbf{X}\right)_{\mathcal{X}}=-(\mathbf{u}, \mathbf{v})_{H}-\partial_{x} u(L) v(L) \tag{29}
\end{equation*}
$$

By definition of $\mathcal{D}\left(\mathcal{A}_{g}\right)$ and Assumption 2, $-\partial_{x} u(L) v(L)=$ $\left[g\left(v_{1}(L)\right)-g\left(v_{2}(L)\right)\right] v(L) \geq 0$, hence

$$
\begin{equation*}
\left(\mathcal{A}_{g}\left(\mathbf{X}_{1}\right)-\mathcal{A}_{g}\left(\mathbf{X}_{2}\right), \mathbf{X}\right)_{\mathcal{X}} \geq-(\mathbf{u}, \mathbf{v})_{H} \tag{30}
\end{equation*}
$$

From (29) we then obtain

$$
\begin{align*}
\left(\mathcal{A}_{g}\left(\mathbf{X}_{1}\right)-\right. & \left.\mathcal{A}_{g}\left(\mathbf{X}_{2}\right)+\mathbf{X}, \mathbf{X}\right)_{\mathcal{X}} \\
& \geq-(\mathbf{u}, \mathbf{v})_{H}+\|\mathbf{u}\|_{V}^{2}+\|\mathbf{v}\|_{H}^{2}  \tag{31}\\
& \geq-(\mathbf{u}, \mathbf{v})_{H}+\frac{1}{2}\|\mathbf{u}\|_{H}^{2}+\frac{1}{2}\|\mathbf{v}\|_{H}^{2} \geq 0
\end{align*}
$$

which is the desired result.
Step 2: Surjectivity. Take $\mathbf{Y}=\left[\mathbf{f}_{1}, \mathbf{f}_{2}\right] \in \mathcal{X}$. Let us prove that there exists $\mathbf{X}=[\mathbf{u}, \mathbf{v}] \in \mathcal{D}\left(\mathcal{A}_{g}\right)$ such that $\mathcal{A}_{g}(\mathbf{X})+$ $2 \mathbf{X}=\mathbf{Y}$, i.e.

$$
\left\{\begin{array}{lc}
\left(-\partial_{x x} u,-\partial_{x} u(0)\right)+2 \mathbf{v}=\mathbf{f}_{2} & \text { in } H  \tag{32}\\
2 \mathbf{u}-\mathbf{v}=\mathbf{f}_{1} & \text { in } V
\end{array}\right.
$$

Replacing $\mathbf{u}$ with $\left(\mathbf{f}_{1}+\mathbf{v}\right) / 2$ in (32) and using the condition $\partial_{x} u(L)=-g(v(L))$, we may start by finding $\mathbf{v} \in V$ such that for all $\mathbf{w} \in V$,

$$
\begin{align*}
\frac{1}{2} a(\mathbf{v}, \mathbf{w})+\frac{1}{2} g(v(L)) & w(L)+2(\mathbf{v}, \mathbf{w})_{H}  \tag{33}\\
& =-\frac{1}{2} a\left(\mathbf{f}_{1}, \mathbf{w}\right)+\left(\mathbf{f}_{2}, \mathbf{w}\right)_{H}
\end{align*}
$$

Define $\Theta: V \rightarrow V^{\prime}$ by, for all $\mathbf{v}, \mathbf{w} \in V$,

$$
\begin{equation*}
\langle\Theta(\mathbf{v}), \mathbf{w}\rangle_{V} \triangleq a(\mathbf{v}, \mathbf{w})+g(v(L)) w(L)+4(\mathbf{v}, \mathbf{w})_{H} \tag{34}
\end{equation*}
$$

and also $L \in V^{\prime}$ by $\langle L, \mathbf{w}\rangle_{V} \triangleq a\left(\mathbf{f}_{1}, \mathbf{w}\right)+2\left(\mathbf{f}_{2}, \mathbf{w}\right)_{H}$. Then, reformulating (33), we seek $\mathbf{v} \in V$ such that

$$
\begin{equation*}
\Theta(\mathbf{v})=L \quad \text { in } V^{\prime} \tag{35}
\end{equation*}
$$

We wish to apply Lemma 3 given in Appendix to $\Theta$. As a Hilbert space, $V$ is reflexive; it is also separable - see Lemma 1. For all $\mathbf{w}_{1}, \mathbf{w}_{2} \in V,\left\langle\Theta\left(\mathbf{w}_{1}\right)-\right.$ $\left.\Theta\left(\mathbf{w}_{2}\right), \mathbf{w}_{1}-\mathbf{w}_{2}\right\rangle_{V}=a\left(\mathbf{w}_{1}-\mathbf{w}_{2}, \mathbf{w}_{1}-\mathbf{w}_{2}\right)+\left[g\left(w_{1}(L)\right)-\right.$ $\left.g\left(w_{2}(L)\right)\right]\left(w_{1}(L)-w_{2}(L)\right)+4\left\|\mathbf{w}_{1}-\mathbf{w}_{2}\right\|_{H}^{2} \geq 0$, which proves the monotonicity. Items (2) and (3) are easily verified using continuity arguments. Finally, let $\mathbf{w} \in V$; we have $\langle L, \mathbf{w}\rangle_{V^{\prime}, V} \leq\|L\|_{V^{\prime}}\|\mathbf{w}\|_{V}$ and $\langle\Theta(\mathbf{w}), \mathbf{w}\rangle_{V^{\prime}, V} \geq\|\mathbf{w}\|_{V}^{2}$, so that $\langle\Theta(\mathbf{w}), \mathbf{w}\rangle_{V^{\prime}, V}-\langle L, \mathbf{w}\rangle_{V^{\prime}, V} \rightarrow+\infty$ when $\|\mathbf{w}\|_{V} \rightarrow$ $+\infty$, which implies the desired property. We deduce from Lemma 3 that there exists $\mathbf{v} \in V$ such that (35) holds. Then, let $\mathbf{u} \triangleq\left(\mathbf{f}_{1}+\mathbf{v}\right) / 2 \in V$. For all $\mathbf{w} \in V$, we have

$$
\begin{equation*}
a(\mathbf{u}, \mathbf{w})+g(v(L)) w(L)+2(\mathbf{v}, \mathbf{w})_{H}=\left(\mathbf{f}_{2}, \mathbf{w}\right)_{H} \tag{36}
\end{equation*}
$$

Since $f_{2}$ belongs to $L^{2}(\Omega)$, taking $\phi=(\phi, 0)$ in (36), where $\phi$ is an arbitrary element of $\mathcal{C}_{c}^{\infty}(\Omega)$, allows us to prove that $u \in H^{2}(\Omega)$. Now, integrating by parts in (36) gives

$$
\begin{align*}
\partial_{x} u(L) w(L)-\partial_{x} u(0) w(0)+g(v(L)) & w(L)+2 v(0) w(0) \\
& =f_{2}(0) w(0) \tag{37}
\end{align*}
$$

holding for all $\mathbf{w} \in V$. Pick a triangular function $\rho$ such that $\rho(0)=1$ and $\rho(L)=0$. Then $(\rho, 1) \in V$. Evaluating (37) with $\rho(x)$ and then $\rho(L-x)$ yields $-\partial_{x} u(0)+2 v(0)=f_{2}(0)$ and $\partial_{x} u(L)=-g(v(L))$, hence $\mathbf{X}=(\mathbf{u}, \mathbf{v}) \in \mathcal{D}\left(\mathcal{A}_{g}\right)$ and $\mathcal{A}_{g}(\mathbf{X})+2 \mathbf{X}=\mathbf{Y}$, which concludes the proof.

In order to consider any initial datum in the state space $\mathcal{X}$, we complete Proposition 2 with the following lemma.
Lemma 2. $\mathcal{D}\left(\mathcal{A}_{g}\right)$ is dense in $\mathcal{X}$.
We can finally prove the Hadamard theorem. Considering the evolution equation 27 with a sequence of approximate initial data in $\mathcal{D}\left(\mathcal{A}_{g}\right)$ provides strong solutions that converge in $\mathcal{C}([0, T], \mathcal{X})$. Then, one has to prove that the limit is a weak solution. For that purpose, using a multiplier method,
we retrieve a standard hidden regularity property of the wave equation.

Proof of Theorem 1 . We split the proof into four steps.
Step 1: Approximate solutions. Since $\mathcal{D}\left(\mathcal{A}_{g}\right)$ is dense in $\mathcal{X}$, we can pick a sequence of vectors $\mathbf{X}_{n}^{0}=\left[\mathbf{u}_{n}^{0}, \mathbf{u}_{n}^{1}\right] \in$ $\mathcal{D}\left(\mathcal{A}_{g}\right)$ that converges to $\left[\mathbf{u}^{0}, \mathbf{u}^{1}\right]$ in $\mathcal{X}$. For each $n \geq 0$, there exists a (unique) strong solution $\mathbf{X}_{n} \in W^{1, \infty}\left(\mathbb{R}^{+}, \mathcal{X}\right)$ to the following reformulation of (27):

$$
\left\{\begin{array}{l}
\dot{\mathbf{X}}(t)+\mathcal{A}_{g}(\mathbf{X}(t))+\mathbf{X}(t)=\mathcal{F}(\mathbf{X}(t))+\mathbf{X}(t)  \tag{38}\\
\mathbf{X}(0)=\mathbf{X}_{n}^{0}
\end{array}\right.
$$

Indeed, Proposition 2 states that $\mathcal{A}_{g}+\mathrm{id}$ is maximal monotone, and $\mathcal{F}+$ id is still Lipschitz. Existence is given by [11, Corollary 4.1]. For all $n \geq 0, \mathbf{X}_{n}=\left[\mathbf{u}_{n}, \mathbf{u}_{n}^{\prime}\right]$ verifies $~^{4}$

$$
\begin{equation*}
\mathbf{u}_{n}^{\prime \prime}(t)-\left(\partial_{x x} u_{n}(t), \partial_{x} u_{n}(0, t)\right)=B\left(\mathbf{u}_{n}^{\prime}(t)\right) \quad \text { in } H \tag{39}
\end{equation*}
$$

in the sense of strong differentiation, for a.e. $t>0$, with $\mathbf{X}_{n}$ taking values in $\mathcal{D}\left(\mathcal{A}_{g}\right)$. We pick $\mathbf{w}$ in $V$ and take the scalar product of (40) with $w$. With an integration by parts, and using the definition of $\mathcal{D}\left(\mathcal{A}_{g}\right)$, we obtain

$$
\begin{array}{r}
\left(\mathbf{u}_{n}^{\prime \prime}(t), \mathbf{w}\right)_{H}+a\left(\mathbf{u}_{n}(t), \mathbf{w}\right)=\left(B\left(\mathbf{u}_{n}^{\prime}(t)\right), \mathbf{w}\right)_{H} \\
-g\left(\partial_{t} u_{n}(L, t)\right) w(L) \tag{40}
\end{array}
$$

for a.e. $t>0$, which means that, $\mathbf{w}$ being arbitrary,

$$
\begin{equation*}
\mathbf{u}_{n}^{\prime \prime}(t)+A \mathbf{u}_{n}(t)=B\left(\mathbf{u}_{n}^{\prime}(t)\right)-g\left(\partial_{t} u_{n}(L, t)\right) \delta_{L} \quad \text { in } V^{\prime} \tag{41}
\end{equation*}
$$

for a.e. $t>0$. Let $T>0$. Since $\mathbf{u}_{n} \in L^{\infty}(0, T ; V) \cap$ $W^{1, \infty}(0, T ; H)$ and $\mathbf{u}_{n}^{\prime} \in W^{1,2}\left(0, T ; V^{\prime}\right)$, a standard timeregularization argument - omitted here - allows us to infer from 41 that the distributional identity $(8)$ holds for all $\phi \in$ $\mathcal{C}([0, T], V) \cap \mathcal{C}^{1}([0, T], H)$, meaning that $\mathbf{u}_{n}$ is indeed a weak solution to (1) in the sense of Definition 11 with initial data $\left[\mathbf{u}_{n}^{0}, \mathbf{v}_{n}^{0}\right]$. Besides, as a consequence of [11, Corollary 4.1], $\mathbf{X}_{n}$ takes values in $\mathcal{D}\left(\mathcal{A}_{g}\right)$ and $\left\|\mathcal{A}_{g}\left(\mathbf{X}_{n}\right)\right\|_{\mathcal{X}}$ is bounded on $[0, T]$, meaning that

$$
\begin{equation*}
\mathbf{u}_{n} \in L^{\infty}(0, T ; W), \quad \mathbf{u}_{n}^{\prime} \in L^{\infty}(0, T ; V) \tag{42}
\end{equation*}
$$

also, $\partial_{t} u_{n}(0, \cdot)$ is absolutely continuous and $\partial_{t t} u_{n}(0, t)-$ $\partial_{x} u_{n}(0, t)=F\left(\partial_{t} u_{n}(0, t)\right)$ Thus, $u_{n}(0, \cdot)$ belongs to $H^{2}(0, T)$ with the desired weak derivatives. It can be inferred from the distributional identity $(8)$ and the additional regularity given by (42) that $u_{n}(x, t)$ belongs to $H^{2}\left(Q_{T}\right)$ and verifies the wave equation on $Q_{T}$.

Step 2: Uniform convergence. The convergence of the sequence $\left\{\mathbf{X}_{n}\right\}$ in $\mathcal{C}([0, T], \mathcal{X})$ for any $T>0$ follows from the standard monotonicity argument. We denote its limit by $\mathbf{X}=\left[\mathbf{u}, \mathbf{u}^{\prime}\right]$. Coming back to the sequence of approximate strong solutions, for all $(m, n)$ in $\mathbb{N}^{2}$, we write $\mathbf{w}_{m n} \triangleq \mathbf{u}_{m}-$ $\mathbf{u}_{n}$. We also write $\mathbf{w}_{m n}^{0} \triangleq \mathbf{w}_{m n}(0)$ and $\mathbf{w}_{m n}^{1} \triangleq \mathbf{w}_{m n}^{\prime}(0)$. Multiplying the wave equation by $x \partial_{x} w_{m n}(x, t) \in H^{1}\left(Q_{\tau}\right)$,

[^4]we obtain, after several integrations by parts, the following boundary estimate:
\[

$$
\begin{array}{r}
\int_{0}^{\tau}\left|g\left(\partial_{t} u_{m}(L, t)\right)-g\left(\partial_{t} u_{n}(L, t)\right)\right|^{2}+\left|\partial_{t} w_{m n}(L, t)\right|^{2} \mathrm{~d} t \\
=\left.\frac{2}{L} \int_{\Omega} x \partial_{t} w_{m n}(x, t) \partial_{x} w_{m n}(x, t) \mathrm{d} x\right|_{0} ^{\tau} \\
\quad+\frac{1}{L} \iint_{Q_{\tau}}\left|\partial_{t} w_{m n}\right|^{2}+\left|\partial_{x} w_{m n}\right|^{2} \tag{43}
\end{array}
$$
\]

From (43), using Cauchy-Schwarz and Young inequalities, and the Lipschitz property of $B$, we obtain the following:

$$
\begin{align*}
\int_{0}^{\tau} \mid g\left(\partial_{t} u_{m}(L, t)\right)- & \left.g\left(\partial_{t} u_{n}(L, t)\right)\right|^{2}+\left|\partial_{t} w_{m n}(L, t)\right|^{2} \mathrm{~d} t \\
& \leq M \sup _{t \in[0, \tau]}\left[\left\|\mathbf{w}_{m n}\right\|_{V}^{2}+\left\|\mathbf{w}_{m n}\right\|_{H}^{2}\right] \tag{44}
\end{align*}
$$

for all $\tau \in[0, T]$ and some $M>0$. As a consequence, $\left\{\partial_{t} u_{n}(L, \cdot)\right\}$ and $\left\{g\left(\partial_{t} u_{n}(L, \cdot)\right)\right\}$ are Cauchy sequences in $L^{2}(0, T)$. Since $u_{n}(L, \cdot)$ converges (in particular) in $L^{2}(0, T)$, it follows that $\left\{u_{n}(L, \cdot)\right\}$ is a Cauchy sequence in $H^{1}(0, T)$, meaning that $u(L, \cdot)$ belongs in fact to $H^{1}(0, T)$, which is the desired hidden regularity property. By continuity of $g$ and unicity of the limit, $g\left(\partial_{t} u_{n}(L, \cdot)\right)$ converges to $g\left(\partial_{t} u(L, \cdot)\right)$ in $L^{2}(0, T)$ - one has to consider a subsequence converging for a.e. $t \in(0, T)$.

Step 3: Existence of weak solutions. Each $\mathbf{u}_{n}$ verifies the distributional identity (8) for all test-functions $\phi$ in $\mathcal{C}_{c}(0, T ; V) \cap \mathcal{C}_{c}^{1}(0, T ; H)$. We pick $\phi$ and we let $n \rightarrow+\infty$. As $\left[\mathbf{u}_{n}, \mathbf{u}_{n}^{\prime}\right]$ converges to $\left[\mathbf{u}, \mathbf{u}^{\prime}\right]$ in $\mathcal{C}([0, T], \mathcal{X}), \partial_{t} u_{n}(L, \cdot)$ converges to $\partial_{t} u(L, \cdot)$ in $L^{2}(0, T)$, and $g\left(\partial_{t} u_{n}(L, \cdot)\right)$ converges to $g\left(\partial_{t} u(L, \cdot)\right)$ in $L^{2}(0, T)$, we obtain that $\mathbf{u}$ verifies (8) as well and is therefore a weak solution to (1).

Step 4: Unicity and Hadamard continuity. Let $w \triangleq$ $\mathbf{u}_{1}-\mathbf{u}_{2}$, where $\mathbf{u}_{1}$ and $\mathbf{u}_{2}$ are two arbitrary weak solutions. The following energy identity holds:

$$
\begin{align*}
& \left.\mathcal{E}^{\mathbf{w}}(t)\right|_{0} ^{\tau}=\int_{0}^{\tau}\left(B\left(\mathbf{u}_{1}^{\prime}(t)\right)-B\left(\mathbf{u}_{2}^{\prime}(t)\right), \mathbf{w}^{\prime}(t)\right)_{H}  \tag{45}\\
& -\left[g\left(\partial_{t} u_{1}(L, t)\right)-g\left(\partial_{t} u_{2}(L, t)\right)\right] \partial_{t} w(L, t) \mathrm{d} t
\end{align*}
$$

for any $0 \leq \tau$. It is obtained by multipling the wave equation by $\partial_{t} w(x, t)$ and integrating over $Q_{\tau}$. From there, using the fact that $g$ is nondecreasing and $B$ is Lipschitz, since we also have $\left.\|\mathbf{w}(t)\|_{H}^{2}\right|_{0} ^{\tau} \leq 2 \int_{0}^{\tau}\|\mathbf{w}\|_{H}\left\|\mathbf{w}^{\prime}\right\|_{H}$, we can apply the standard Grönwall argument to obtain unicity and continuity with respect to initial conditions, proving that the operators $\mathcal{S}_{t}$ introduced in (9) are well-defined and continuous. At this point, Proposition 1 is also proved.

## VI. Conclusion and perspectives

In this paper, the well-posedness of a wave equation with a nonlinear dynamic boundary condition and nonlinear Neuman-type feedback is proved, in a variational framework. Using a proportional feedback law, exponential stabilization is achieved under suitable assumptions.

In the linear case, an infinite-dimensional frequencydomain approach would be interesting to tackle the stabilization problem and obtain sharp conditions. Besides, adding an integral action to the feedback law should be considered. Also, the experimental device of [8] simulating drilling dynamics may be used to test the model in real experiments.

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## Appendix

We use the following existence result, which is an application of the Galerkin method and the Brouwer fixedpoint theorem in finite dimension - see [11, Lemma 2.1 and Theorem 2.1] for the proof.

Lemma 3. Let $E$ be a separable reflexive Banach space and $f \in E^{\prime}$. Assume $\Theta: E \rightarrow E^{\prime}$ verifies the following assumptions:

1) $\Theta$ is monotone i.e. $\left\langle\Theta\left(x_{1}\right)-\Theta\left(x_{2}\right), x_{1}-x_{2}\right\rangle_{E} \geq 0$ for all $x_{1}, x_{2} \in E$;
2) $\Theta$ is bounded i.e. $S \subset E$ bounded implies $\Theta(S)$ bounded in $E^{\prime}$;
3) For all $x_{1}, x_{2} \in E$, the scalar function $t \mapsto\left\langle\Theta\left(x_{1}+\right.\right.$ $\left.\left.t x_{2}\right), x_{2}\right\rangle_{E^{\prime}, E}$ is continuous.
If for some $\rho>0,\|x\|_{E}>\rho$ implies $\langle\Theta(x), x\rangle_{E^{\prime}, E}>$ $\langle f, x\rangle_{E^{\prime}, E}$, then there exists $x \in E$ such that $\Theta(x)=f$.

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[^1]:    ${ }^{1}$ In the language of systems theory, one may say that the Neuman input operator is unbounded with respect to the state space $\mathcal{X}$.

[^2]:    ${ }^{2}$ It is a specificity of the one-dimensional case that no additional growth assumption on $g$ is required to obtain weak solutions for such feedback.

[^3]:    ${ }^{3}$ Alternatively, if $F$ has a stabilizing effect, e.g. $F(0)=0, F$ nonincreasing, it can be relaxed to $F$ locally Lipschitz - with appropriate modifications in the proof of well-posedness.

[^4]:    ${ }^{4}$ One should verify that the second coordinate is indeed the weak derivative of the first coordinate.

