SENSITIVITY ANALYSIS FOR MPC AND AN EXPONENTIAL TURNPIKE THEOREM FOR LINEAR QUADRATIC OPTIMAL CONTROL OF GENERAL EVOLUTION EQUATIONS *

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Abstract. We analyze the sensitivity of linear quadratic optimal control problems governed by general evolution equations with bounded or admissible control operator. We show, that if the problem is stabilizable and detectable, the solution of the extremal equation can be bounded by the right-hand side including initial data with the bound being independent of the time horizon. Consequently, the influence of perturbations of the extremal equations decays exponentially in time. This property can for example be used to construct efficient space and time discretizations for a Model Predictive Control scheme. Furthermore, a turnpike property for unbounded but admissible control of general semigroups can be deduced.

Key words. Sensitivity Analysis, Turnpike Property, Model Predictive Control

AMS subject classifications. 49K20, 49K40, 93D20

1. Introduction. In this paper we consider optimal control of linear evolution equations in infinite dimensional spaces. It seems intuitively clear that perturbations of the data that occur in the far future affect the optimal control at present time only slightly. This might serve as an explanation, why closed-loop control techniques, which only take present information into account, often work in a nearly optimal way, even in the long run. The aim of this work is to give a rigorous and quantitative justification to this statement.

More specifically, we analyze the sensitivity of the solutions to linear quadratic optimal control problems, characterized by the first order optimality conditions with respect to the problem data, i.e., source terms and initial data of the state and adjoint equation. In the first part of our paper we are concerned with the question on how the behavior of perturbations of the source terms over time influences the solutions' behavior over time. It turns out that if the dynamics are stabilizable and detectable, the perturbations' influence decays exponentially as the distance in time grows.

Our paper generalizes recent results in this direction from parabolic equations [15], which enjoy very favorable smoothing properties, to general evolution equations with much less structure, which in particular comprises hyperbolic equations. This makes the analysis much more delicate. For example, in [15], the boundary control or boundary observation case could be dealt with in a straightforward fashion, as the state and adjoint belong to a Sobolev space for a.a. times on which boundary control or observation operators are bounded. This is no longer the case when moving to general evolution equations. Instead the concept of unbounded, but admissible control and observation operators and well-posed linear systems as described e.g. in [31] and [28] has to be employed to cover the case of boundary control. Many general results, such as Riccati theory, are not available in their full strength within this broader scope.

The analysis in this work is based on analyzing the extremal equations of the optimal control problem, avoiding the use of Riccati equations. Compared to [15] we use weaker assumptions and thus obtain estimates in weaker norms, concerning spatial regularity, which is to be expected when going from parabolic to e.g. hyperbolic equations. Concerning temporal regularity, we show uniform

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estimates as well as L_2 -type estimates for perturbations in L_2 and in L_1 .

One important motivation for our sensitivity analysis is Model Predictive Control (MPC). MPC is a control technique, where the solution of an optimal control problem on an indefinite or infinite horizon is split into the successive solution of problems with short time horizons T. However, only a first part up to a time τ , where $\tau \ll T$ is implemented and the process is repeated, with the resulting state as initial condition. Consequently, only the first part of the control has to be computed accurately. The fact that perturbations decay exponentially in time allows to discretize and compute solutions with reduced accuracy towards the end of the time horizon, as discretization errors are perturbations of the extremal equations. We refer to [14] for an in-depth introduction to MPC and to [12] for an overview of its approximation properties.

As a second application of our technique we show an exponential turnpike property. This property is important to understand the behavior of solutions to optimal control problems on large time horizons, which remain close to a steady state solution, the so called turnpike, for the majority of the time. There is a large literature on turnpike behavior, cf. [1, 8, 18, 35, 36, 13, 10, 17, 16]. A particular kind of turnpike behavior is the so called exponential turnpike property, cf. [30, 29, 23, 24, 7, 5]. The proofs establishing the turnpike property in the recent works [5, 29, 15] are based on a stabilizability and detectability assumption on the system. Turnpike theorems in Hilbert spaces were given in [29] for general semigroups with bounded control and observation operators, and for boundary controlled parabolic equations. The proofs in [29], however, make use of the Algebraic Riccati Equation, a theory, which is well established for admissible boundary control of parabolic equations, but not for general evolution equations. Here, we show a turnpike result for unbounded but admissible control of non-parabolic equations which was not available up to now. This is possible, since we avoid using Riccati theory in our approach. Our analysis also sheds light on the close connection of exponential sensitivity analysis and the turnpike property, both emerging from the boundedness of the operator corresponding to the extremal equations. This becomes clear by comparing the abstract scaling results in Theorem 3.2 and Theorem 6.2.

This paper is structured as follows. First, in Section 2 we present our theoretical framework, the optimal control problem and its optimality conditions. In Section 3, we derive a general result on the error propagation over time in Theorem 3.2, under the assumption that various norms of the extremal equations' solution operator, which itself may indeed depend on the horizon T, can be bounded independently of T. Then, in Section 4 we show that such a T-independent bound on these operator norms holds, if the dynamics are exponentially stabilizable and detectable. This assumption enters into the construction of special test functions, similar to [23] and [29], that decay exponentially. In Section 5, we will extend the results to unbounded but admissible control operators and discuss the necessary modifications to the proofs. In Section 6 we show our turnpike result Theorem 6.2 for general evolution equations. Here, too, we allow for unbounded but admissible control for general C_0 -semigroups. Then, in Section 7, under an exact controllability assumption, we extend our results to the case of a terminal condition on the state. Last, we present two examples that fulfill the assumptions of our analysis, namely the interior control of a heat equation, and the Dirichlet boundary control of a wave equation.

2. Setting and preliminaries. We consider linear quadratic optimal control problems in a Hilbert space X of the following type.

Problem 2.1.

(2.1)
$$\min_{(y,u)} \frac{1}{2} \int_0^T ||C(y(t) - y_d)||_Y^2 + ||R(u(t) - u_d)||_U^2 dt$$
$$y' - \Lambda y - Bu = f$$
$$y(0) = y_0,$$

where Y and U are Hilbert spaces for the output and the control and T > 0 is the optimization horizon. The dynamics are given by an evolution equation where $\Lambda : D(\Lambda) \subset X \to X$ is a (possibly unbounded) generator of a C_0 -semigroup $(\mathcal{T}(t))_{t\geq 0}$ on X, $f \in L_1(0,T;X)$ is a source term and $y_0 \in X$ is the initial datum. Moreover, $R \in L(U,U)$ is an elliptic operator, leading to convexity of the problem. Finally, the control and observation operators are bounded, i.e., $B \in L(U,X)$ and $C \in L(X,Y)$.

Only in the first part of this paper we assume that B and C are bounded, in order not to hide the main steps of our proof behind technical details. In Section 5 we will generalize the result to unbounded but admissible control operators.

Whenever we consider a solution to an evolution equation, we refer to the mild solution $y \in C(0,T;X)$, which for (2.1) can be defined by

$$y(t) = \mathcal{T}(t)y_0 + \int_0^t \mathcal{T}(t-s)(Bu(s) + f(s)) ds.$$

We refer to [22, Section 4.2] for an introduction to inhomogeneous abstract Cauchy problems. In the following, $\|\cdot\|$ resp. $\langle\cdot,\cdot\rangle$ denote the norm resp. the scalar product in the Hilbert spaces X,U and Y. The space C(0,T;X) with norm $\|x\|_{C(0,T;X)}:=\max_{t\in[0,T]}\|x(t)\|$ will be denoted by C(X) and $L_p(0,T;X)$ with norm $\|x\|_{L_p(0,T;X)}^p:=\sum_{t=0}^T\|x(t)\|^pdt$ will be denoted by $L_p(X)$ for $p\in\{1,2\}$ and $L_\infty(0,T;X)$ with norm $\|x\|_{L_\infty(0,T;X)}:=\operatorname{ess\,sup}_{t\in[0,T]}\|x(t)\|$ by $L_\infty(X)$. We will assume that there exists at least one minimizer of Problem 2.1. For the existence of minimizers, we refer to [20, Chapter 3]. Let $(y,u)\in C(X)\times L_2(U)$ be a minimizer of Problem 2.1. Then, by the Pontryagin Maximum Principle, see [20, Chapter 4], there exists $\lambda\in C(X)$ such that

(2.2)
$$C^*Cy - \lambda' - \Lambda^*\lambda = C^*Cy_d$$
$$R^*Ru - B^*\lambda = R^*Ru_d$$
$$y' - \Lambda y - Bu = f$$

in a mild sense along [0,T] with initial resp. terminal conditions $y(0)=y_0$ and $\lambda(T)=0$. We denote by Λ^* the adjoint operator associated with Λ with domain $D(\Lambda^*)$. Defining $Q:=R^*R$ and eliminating the control via the second equation with $u=Q^{-1}B^*\lambda+u_d$ leads to the linear system of equations

(2.3)
$$\underbrace{\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix}}_{-:M} \begin{pmatrix} y \\ \lambda \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ 0 \\ Bu_d + f \\ y_0 \end{pmatrix},$$

where $E_0y := y(0)$ and $E_T\lambda := \lambda(T)$. The operator M corresponds to the two abstract inhomogeneous evolution equations, i.e., the adjoint equation in the first two rows and the state equation in the last two rows and allows for a brief notation of those. The solution operator of this system, which we denote by M^{-1} , maps initial values and source terms for the state and the adjoint equation to the solution. This mapping is well defined as a mapping from $(L_1(X) \times X)^2$ to $C(X)^2$, as for right-hand sides with this regularity, the state and adjoint equations have a unique solution in C(X) each. A central question in the following will be the dependence of this inverse operator's norm on the time T.

3. A sensitivity result. In this section, we present a first sensitivity result which extends the result of [15, Theorem 3.1]. For this, we will refer to the solution $(y, \lambda) \in C(X)^2$ of (2.3) as the exact solution. We will now assume that there is a second pair of variables $(y_p, \lambda_p) \in C(X)^2$ that satisfies the perturbed system

$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} y_p \\ \lambda_p \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ 0 \\ Bu_d + f \\ y_0 \end{pmatrix} + \begin{pmatrix} \varepsilon_1 \\ 0 \\ \varepsilon_2 \\ 0 \end{pmatrix}$$

for $(\varepsilon_1, \varepsilon_2) \in L_1(X)^2$. This solution will be referred to as the *perturbed solution*. The terms ε_1 and ε_2 are perturbations of the dynamics which could stem from discretization errors in time or space. In this subsection we will give an estimate for the norm of the difference of (y_p, λ_p) and (y, λ) . It follows by linearity that the difference between exact and perturbed solution $(\delta y, \delta \lambda) := (y_p - y, \lambda_p - \lambda)$ solves

(3.1)
$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} \delta y \\ \delta \lambda \end{pmatrix} = \begin{pmatrix} \varepsilon_1 \\ 0 \\ \varepsilon_2 \\ 0 \end{pmatrix}.$$

The question we want to answer is the following. How does the behavior of the perturbations ε_1 and ε_2 over time influence the behavior of δy and $\delta \lambda$? Theorem 3.2 gives an answer towards that direction. We will prove a scaling result for evolution equations first. It is a well-known result that, if a semigroup $(\mathcal{T}(t))_{t\geq 0}$ has generator Λ , the scaled semigroup $(e^{-\mu t}\mathcal{T}(t))_{t\geq 0}$ has generator $\Lambda - \mu I$ [9, p.60] with the same domain as Λ , as the domain does not change under bounded perturbations, cf. [9, Chapter III]. In our context, we want to put emphasis on the impact of scaling on the source terms ε_1 and ε_2 .

Lemma 3.1. Assume $\delta y, \delta \lambda \in C(X)$ solve the abstract Cauchy problems

$$\delta y' - \Lambda \delta y = \varepsilon_1$$
$$-\delta \lambda' - \Lambda^* \delta \lambda = \varepsilon_2$$

with initial resp. terminal conditions $\delta y(0) = \delta y_0$, $\delta \lambda(T) = \delta \lambda_T$. Then $\widetilde{\delta y}(t) = e^{-\mu t} \delta y(t)$ and $\widetilde{\delta \lambda}(t) = e^{-\mu t} \delta \lambda(t)$ solve

(3.2)
$$\widetilde{\delta y}' - (\Lambda - \mu I)\widetilde{\delta y} = e^{-\mu t} \varepsilon_1 \\ -\widetilde{\delta \lambda}' - (\Lambda + \mu I)^* \widetilde{\delta \lambda} = e^{-\mu t} \varepsilon_2$$

with $\widetilde{\delta y}(0) = \delta y_0$ and $\widetilde{\delta \lambda}(T) = e^{-\mu T} \delta \lambda_T$.

Proof. Using the equivalence of the concept of weak and mild solution, see [20, Section 5.1] and [2], we will use the weak formulation in this case, which is for $\varphi \in D(\Lambda^*)$ given by

$$\frac{d}{dt}\langle \delta y(t), \varphi \rangle - \langle \delta y(t), \Lambda^* \varphi \rangle = \langle \varepsilon_1(t), \varphi \rangle$$

for a.e. $t \in [0,T]$. Inserting $\delta y = e^{\mu t} \widetilde{\delta y}$ and applying the chain rule yields the result for the state equation. The equation for the scaled adjoint variable follows analogously.

The following theorem gives a preliminary result on the propagation of perturbations. It gives an estimate on the scaled difference of exact and perturbed solution under the assumption that different solution operator norms can be bounded independently of T.

THEOREM 3.2. Assume $(\delta y, \delta \lambda) \in C(X)^2$ solves (3.1), where $\varepsilon_1, \varepsilon_2 \in L_1(X)$. Assume the solution operator's norms

(3.3)
$$\|M^{-1}\|_{(L_1(X)\times X)^2\to C(X)^2}, \qquad \|M^{-1}\|_{(L_2(X)\times X)^2\to C(X)^2}, \\ \|M^{-1}\|_{(L_1(X)\times X)^2\to L_2(X)^2}, \qquad \|M^{-1}\|_{(L_2(X)\times X)^2\to L_2(X)^2}$$

can be bounded independently of T. Then there is a constant $c \ge 0$ and a scaling factor $0 < \mu$, both independent of T, such that defining

$$\rho := \|e^{-\mu t} \varepsilon_1(t)\|_E + \|e^{-\mu t} \varepsilon_2(t)\|_E$$

for $E := L_1(X)$ or $E := L_2(X)$, we have

(3.4)
$$||e^{-\mu t} \delta y||_{L_2(X)} + ||e^{-\mu t} \delta \lambda||_{L_2(X)} \le c\rho,$$

$$||e^{-\mu t} \delta u||_{L_2(U)} \le c\rho$$

and

(3.5)
$$||e^{-\mu t} \delta y||_{C(X)} + ||e^{-\mu t} \delta \lambda||_{C(X)} \le c\rho, \\ ||e^{-\mu t} \delta u||_{L_{\infty}(U)} \le c\rho.$$

Proof. For $\mu > 0$ we introduce scaled variables $\widetilde{\delta y}(t) := e^{-\mu t} \delta y(t)$, $\widetilde{\delta \lambda}(t) := e^{-\mu t} \delta \lambda(t)$, $\widetilde{\varepsilon}_1(t) := e^{-\mu t} \varepsilon_1(t)$ and $\widetilde{\varepsilon}_2(t) := e^{-\mu t} \varepsilon_2(t)$ and apply Lemma 3.1. This yields

$$(3.1) \qquad \Longleftrightarrow \qquad \left(\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} + \mu \underbrace{\begin{pmatrix} 0 & -I \\ 0 & 0 \\ I & 0 \\ 0 & 0 \end{pmatrix}}_{=:P} \right) \begin{pmatrix} \widetilde{\delta y} \\ \widetilde{\delta \lambda} \end{pmatrix} = \begin{pmatrix} \widetilde{\varepsilon}_1 \\ 0 \\ \widetilde{\varepsilon}_2 \\ 0 \end{pmatrix}.$$

Introducing $\widetilde{z}:=(\widetilde{\delta y},\widetilde{\delta\lambda}),\,\widetilde{\varepsilon}:=(\widetilde{\varepsilon}_1,0,\widetilde{\varepsilon}_2,0)$ we compute formally

$$(3.6) \hspace{1cm} (M+\mu P)\widetilde{z}=\widetilde{\varepsilon} \hspace{0.3cm} \Rightarrow \hspace{0.3cm} (I+\mu PM^{-1})M\widetilde{z}=\widetilde{\varepsilon} \hspace{0.3cm} \Rightarrow \hspace{0.3cm} \widetilde{z}=M^{-1}(I+\mu PM^{-1})^{-1}\widetilde{\varepsilon}.$$

Next we expand $(I + \mu P M^{-1})^{-1}$ into a Neumann series, cf. [19, Theorem 2.14]. In the following denote $W = (L_2(X) \times X)^2$. We choose $\mu > 0$, such that $\beta := \mu \|M^{-1}\|_{W \to L_2(X)^2} < 1$. By

assumption $||M^{-1}||_{W\to L_2(X)^2}$ is bounded independently of T, so we can choose $\mu>0$ independently of T.

Since $||P||_{L_2(X)^2 \to W} \le 1$, it follows $||\mu P M^{-1}||_{W \to W} \le \beta < 1$. Neumann series expansion of $(I + \mu P M^{-1})^{-1}$ yields

(3.7)
$$||(I + \mu P M^{-1})^{-1}||_{W \to W} \le \sum_{i=0}^{\infty} ||(\mu P M^{-1})^i||_{W \to W} \le \sum_{i=0}^{\infty} \beta^i = \frac{1}{1 - \beta}.$$

This argumentation already yields the desired result in the L_2 case:

$$\|\tilde{z}\|_{L_2(X)}^2 \leq \|M^{-1}\|_{W \to L_2(X)^2} \|(I + \mu P M^{-1})^{-1}\|_{W \to W} \|\tilde{\varepsilon}\|_W.$$

To prove the general case we use the following alternative representation, which can be verified by premultiplication with $(I + \mu PM^{-1})$:

$$(I + \mu P M^{-1})^{-1} = I - (I + \mu P M^{-1})^{-1} \mu P M^{-1}.$$

Now we estimate the operator norm for $S = C(X)^2$ or $S = L_2(X)^2$ via

$$||M^{-1}(I + \mu P M^{-1})^{-1}||_{(E \times X)^2 \to S} = ||M^{-1} - M^{-1}(I + \mu P M^{-1})^{-1} \mu P M^{-1}||_{(E \times X)^2 \to S}$$

$$\leq ||M^{-1}||_{(E \times X)^2 \to S} + ||M^{-1}||_{W \to S} ||(I + \mu P M^{-1})^{-1}||_{W \to W} ||\mu P M^{-1}||_{(E \times X)^2 \to W}$$

$$\leq ||M^{-1}||_{(E \times X)^2 \to S} + \frac{\mu ||M^{-1}||_{W \to S} ||M^{-1}||_{(E \times X)^2 \to L_2(X)^2}}{1 - \beta}$$

using $||P||_{L_2(X)^2 \to W} \leq 1$ and (3.7). Thus, it follows by (3.6) with taking norms, that

$$(3.8) \quad \begin{aligned} \|\widetilde{z}\|_{L_{2}(X)^{2}} &\leq \left(\|M^{-1}\|_{(E \times X)^{2} \to L_{2}(X)^{2}} + \frac{\mu \|M^{-1}\|_{W \to L_{2}(X)^{2}} \|M^{-1}\|_{(E \times X)^{2} \to L_{2}(X)^{2}}}{1 - \beta} \right) \|\widetilde{\varepsilon}\|_{E^{2}}, \\ \|\widetilde{z}\|_{C(X)^{2}} &\leq \left(\|M^{-1}\|_{(E \times X)^{2} \to C(X)^{2}} + \frac{\mu \|M^{-1}\|_{W \to C(X)^{2}} \|M^{-1}\|_{(E \times X)^{2} \to L_{2}(X)^{2}}}{1 - \beta} \right) \|\widetilde{\varepsilon}\|_{E^{2}}. \end{aligned}$$

Using the assumption on the T-independent bound on the operator norms and the definition $\rho := \|\widetilde{\varepsilon}\|_{E^2}$, we obtain the result for the state and the adjoint by going back to the original variables via $\widetilde{z} = (e^{-\mu t}\delta y, e^{-\mu t}\delta \lambda)$. For the control we conclude

(3.9)
$$\|e^{-\mu t}\widetilde{\delta u}\|_{L_2(U)} = \|e^{-\mu t}Q^{-1}B^*\widetilde{\delta \lambda}\|_{L_2(U)} \le \|Q^{-1}B^*\|_{X\to U} \|\|\widetilde{\delta \lambda}\|_{L_2(X)} \le c\rho,$$

$$\|e^{-\mu t}\widetilde{\delta u}(t)\| = \|e^{-\mu t}Q^{-1}B^*\widetilde{\delta \lambda}(t)\| \le \|Q^{-1}B^*\|_{X\to U} \|\delta \lambda(t)\| \le c\rho$$

for a.a. $t \in [0, T]$, where we used the bound on $||e^{-\mu t}\delta\lambda||_{L_2(X)}$ and $||e^{-\mu t}\delta\lambda||_{C(X)}$ and the fact that B and Q do not involve time derivatives. This yields (3.4) and (3.5).

Remark 3.3. We will briefly comment on the Neumann series occurring in (3.7). The operator $(I+\mu PM^{-1})^{-1}$ can be represented by its Neumann series, i.e., $(I+\mu PM^{-1})^{-1} = \sum_{i=0}^{\infty} (-\mu PM^{-1})^k$, see [19, Theorem 2.14]. We provide an illustration for the summand for k=2, i.e., $(\mu PM^{-1})^2 = \mu PM^{-1}\mu PM^{-1}$. The application of this operator can be interpreted as the following. M^{-1} solves the corresponding Cauchy problems with right-hand side including initial and terminal condition. Afterwards, the operator μP maps the solutions to source terms scaled by μ , i.e., $\mu P(\delta y, \delta \lambda) = (-\mu \delta \lambda, 0, \mu \delta y, 0)$. This right-hand side then enters M^{-1} again, the Cauchy problems are solved with zero initial data and source terms $-\mu \delta \lambda$ and $\mu \delta y$ and the process is repeated.

The crucial assumption in Theorem 3.2 is that the operator norms in (3.3) can be bounded independently of T. In Section 4, we will show, that this is indeed the case, if the dynamics are exponentially stabilizable and detectable. It will turn out that all bounds in (3.3) can be shown simultaneously.

Before we head to the next part, we recall an integration by parts result for solutions expressed via semigroups.

Lemma 3.4. Consider the abstract Cauchy problems

$$x'_1 = \Lambda x_1 + f_1,$$
 $x_1(0) = x_{10},$
 $-x'_2 = \Lambda^* x_2 + f_2,$ $x_2(0) = x_{20},$

where $x_{10}, x_{20} \in X$ and $f_1, f_2 \in L_1(X)$. $\Lambda : D(\Lambda) \to X$ is the generator of a C_0 -semigroup $(\mathcal{T}(t))_{t \geq 0}$ on X. Consider the associated mild solutions $x_1, x_2 \in C(0, T; X)$ given by

$$x_1(t) = \mathcal{T}(t)x_{10} + \int_0^t \mathcal{T}(t-s)f_1(s) \, ds \qquad \forall t \in [0,T],$$

$$x_2(t) = \mathcal{T}^*(T-t)x_{20} + \int_t^T \mathcal{T}^*(s-t)f_2(s) \, ds \qquad \forall t \in [0,T].$$

Then, for all $0 \le s \le t \le T$,

$$\langle x_1(t), x_2(t) \rangle - \langle x_1(s), x_2(s) \rangle = \int_s^t (\langle x_2, f_1 \rangle - \langle x_1, f_2 \rangle) ds.$$

Proof. For a proof of this property, we refer to [20, Proposition 5.7, p.69].

4. T-independent bounds for the solution operator. In this section we will derive T-independent bounds on the norm of the solution operator M^{-1} , which is a central assumption in the abstract scaling result of Theorem 3.2. Since [0, T] is bounded, we have the continuous embeddings

$$C(0,T;X) \hookrightarrow L_2(0,T;X) \hookrightarrow L_1(0,T;X)$$
.

Hence, we may equip $L_1(0,T;X)$ with the equivalent norm

$$||v||_{1\vee 2} := \min\{||v||_{L_1(X)}, ||v||_{L_2(X)}\}$$

(setting $||v||_{L_2(X)} = \infty$ if $v \notin L_2(0,T;X)$), which satisfies

$$\min\left\{1, \frac{1}{\sqrt{T}}\right\} \|v\|_{L_1(X)} \le \|v\|_{1\vee 2} \le \|v\|_{L_1(X)}.$$

Likewise, we equip C(0,T;X) with the equivalent norm

$$||v||_{2\wedge\infty} := \max\{||v||_{L_2(X)}, ||v||_{C(X)}\},$$

satisfying

$$\|v\|_{C(X)} \leq \|v\|_{2\wedge\infty} \leq \max\left\{1,\sqrt{T}\right\}\|v\|_{C(X)}.$$

Obviously, the equivalence of norms deteriorates for $T \to \infty$. We have the following Hölder-like inequality for these norms:

LEMMA 4.1. Let $v \in C(X)$ and $w \in L_1(X)$. Then,

$$\int_0^T \langle v(s), w(s) \rangle \, ds \le ||v||_{2 \wedge \infty} ||w||_{1 \vee 2}.$$

Proof.

$$\int_{0}^{T} \langle v(s), w(s) \rangle \, ds \le \min\{\|v\|_{C(X)} \|w\|_{L_{1}(X)}, \|v\|_{L_{2}(X)} \|w\|_{L_{2}(X)}\}$$

$$\le \min\{\|v\|_{2 \wedge \infty} \|w\|_{L_{1}(X)}, \|v\|_{2 \wedge \infty} \|w\|_{L_{2}(X)}\} \le \|v\|_{2 \wedge \infty} \|w\|_{1 \vee 2}.$$

The main result of this section will be a T-independent bound for

$$||M^{-1}||_{((L_1(X),||\cdot||_{1\vee 2})\times X)^2\to (C(X),||\cdot||_{2\wedge\infty})^2}$$

This implies all desired T-independent bounds in (3.3), required by Theorem 3.2. To this end, consider mild solutions (y, λ) of the system

$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} y \\ \lambda \end{pmatrix} = \begin{pmatrix} l_1 \\ \lambda_T \\ l_2 \\ y_0 \end{pmatrix}$$

in [0,T], where $l_1, l_2 \in L_1(X)$ and $y_0, \lambda_T \in X$ are given. In short we may write again $z = M^{-1}r$ with $z = (y, \lambda)$ and $r = (l_1, \lambda_T, l_2, y_0)$.

Next, we introduce our main assumption. A semigroup $(\mathcal{T}(t))_{t\geq 0}$ is called exponentially stable, if there exist $M, \mu > 0$, such that $\|\mathcal{T}(t)\|_{L(X,X)} \leq Me^{-\mu t}$ for all t > 0.

Assumption 4.2.

- 1. (Λ, C) is exponentially detectable, i.e., there exists $K_C \in L(Y, X)$ such that the semigroup generated by $\Lambda^* + C^*K_C^*$ is exponentially stable.
- 2. (Λ, B) is exponentially stabilizable, i.e., there exists $K_B \in L(X, U)$ such that the semigroup generated by $\Lambda + BK_B$ is exponentially stable.

The approach in Lemma 4.3 and 4.4 is inspired by the stability estimate in [29, Lemma 2] and [23, Lemma 3.5].

Lemma 4.3. Let $\varphi \in C(0, t; X)$ solve

(4.2)
$$\begin{aligned} -\varphi' &= (\Lambda^* + C^* K_C^*) \varphi \quad on \quad [0, t] \\ \varphi(t) &= y(t) \end{aligned}$$

where K_C^* is a stabilizing feedback for (Λ^*, C^*) . Then, there are constants $M_{\varphi}, k_{\varphi} > 0$ such that for test functions $v \in L_2(0,t;X)$

$$(4.3) \qquad \int_0^t |\langle v(s), \varphi(s) \rangle| \, ds \le \|y(t)\| \frac{M_{\varphi}}{\sqrt{k_{\varphi}}} \sqrt{\int_0^t \|v(s)\|^2 e^{-k_{\varphi}(t-s)} \, ds}.$$

Let $\psi \in C(t,T;X)$ solve

(4.4)
$$\psi' = (\Lambda + BK_B)\psi \quad on [t, T]$$
$$\psi(t) = \lambda(t)$$

where K_B is a stabilizing feedback for (Λ, B) . Then, for test functions $v \in L_2(t, T; X)$

(4.5)
$$\int_{t}^{T} |\langle v(s), \psi(s) \rangle| \, ds \leq \|\lambda(t)\| \frac{M_{\psi}}{\sqrt{k_{\psi}}} \sqrt{\int_{t}^{T} \|v(s)\|^{2} e^{-k_{\psi}(s-t)} \, ds}.$$

Proof. We will first prove (4.3). By exponential stability of the semigroup, there exist $M_{\varphi}, k_{\varphi} > 0$, such that

$$\|\varphi(s)\| \le M_{\varphi} e^{-k_{\varphi}(t-s)} \|y(t)\| \qquad 0 \le s \le t.$$

Using this exponential stability, we get for $v \in L_2(0, t; X)$,

$$\int_0^t |\langle v(s), \varphi(s) \rangle| \, ds \le \int_0^t \|v(s)\| \|\varphi(s)\| \, ds \le \|y(t)\| \int_0^t \|v(s)\| M_{\varphi} e^{-k_{\varphi}(t-s)} \, ds.$$

The last term can be estimated via:

$$\begin{split} \int_0^t \|v(s)\| M_{\varphi} e^{-k_{\varphi}(t-s)} \, ds &= \int_0^t \|v(s)\| M_{\varphi} e^{-\frac{k_{\varphi}}{2}(t-s)} \cdot e^{-\frac{k_{\varphi}}{2}(t-s)} \, ds \\ &\leq \sqrt{\int_0^t \|v(s)\|^2 M_{\varphi}^2 e^{-k_{\varphi}(t-s)} \, ds} \cdot \underbrace{\sqrt{\int_0^t e^{-k_{\varphi}(t-s)} \, ds}}_{<\frac{1}{\sqrt{k_{\varphi}}}}. \end{split}$$

The estimate (4.5) follows analogously.

By using φ and ψ from (4.2) and (4.4), respectively, as test functions for (4.1), we obtain the following pointwise in time identities:

LEMMA 4.4. Let (y, λ) solve (4.1). If φ solves (4.2), then

$$(4.6) \quad ||y(t)||^2 = \int_0^t -\langle K_C C y(s), \varphi(s) \rangle + \langle R^{-*} B^* \lambda(s), R^{-*} B^* \varphi(s) \rangle + \langle l_2(s), \varphi(s) \rangle \, ds + \langle y_0, \varphi(0) \rangle$$

for all $0 \le t \le T$. If ψ solves (4.4), then

$$(4.7) \|\lambda(t)\|^2 = \int_t^T -\langle K_B^* B^* \lambda(s), \psi(s) \rangle - \langle Cy(s), C\psi(s) \rangle + \langle l_1(s), \psi(s) \rangle \, ds + \langle \lambda_T, \psi(T) \rangle$$

for all 0 < t < T.

Proof. We begin with the proof of (4.6). Testing the state equation with $\varphi(s)$ solving (4.2), integration on [0,t] and integration by parts in the sense of Lemma 3.4 on [0,t] with $x_1=y$, $f_1=BQ^{-1}B^*\lambda+l_2, x_2=\varphi, f_2=C^*K_C^*\varphi$ yields

$$\langle y(t), \varphi(t) \rangle - \langle y(0), \varphi(0) \rangle = \int_0^t \langle \varphi(s), BQ^{-1}B^*\lambda(s) + l_2(s) \rangle - \langle y(s), C^*K_C^*\varphi(s) \rangle ds.$$

Rearranging the terms, using the terminal condition $\varphi(t) = y(t)$ and $Q^{-1} = (R^*R)^{-1} = R^{-1}R^{-*}$, we get

$$||y(t)||^2 = \int_0^t -\langle K_C C y(s), \varphi(s) \rangle + \langle R^{-*} B^* \lambda(s), R^{-*} B^* \varphi(s) \rangle + \langle l_2(s), \varphi(s) \rangle \, ds + \langle y_0, \varphi(0) \rangle.$$

The formula (4.7) follows analogous by testing the adjoint equation with ψ solving (4.4).

Based on (4.6) and (4.7) we will derive norm estimates for M^{-1} as a mapping into $L_2(X)^2$ and $C(X)^2$. While the latter turns out to be rather straightforward (see below), the L_2 -estimate requires to integrate (4.6) and (4.7) over [0, T]. The crucial observation is that the integrals on the right hand side of (4.6) and (4.7) can be converted into convolutions with exponentially decaying functions. This will allow us to derive an L_2 -estimate without any constants depending on the time T with the help of the following general lemma:

LEMMA 4.5. For $w \in L_1(0,T;[0,\infty))$, consider

$$h_1(t) := \int_0^t w(s)e^{-k_{\varphi}(t-s)} ds, \text{ where } k_{\varphi} > 0,$$

$$h_2(t) := \int_t^T w(s)e^{-k_{\psi}(s-t)} ds, \text{ where } k_{\psi} > 0.$$

Then, there is a constant $c \geq 0$ independent of T, such that

$$||h_i||_{L_p(0,T)} \le c||w||_{L_1(0,T)}$$
 for $i = 1, 2$ and $1 \le p \le \infty$

Proof. Extending w by 0 from [0,T] to \mathbb{R} and defining $g_1(x) := e^{-k_{\varphi}x}$ for $x \geq 0$ and $g_1(x) = 0$ otherwise, we can write h_1 as a convolution:

$$h_1(t) = (g_1 * w)(t) = \int_{\mathbb{R}} g_1(t-s)w(s) ds,$$

and apply Young's inequality:

$$||h_1||_{L_p(0,T)} = ||g_1 * w||_{L_p(\mathbb{R})} \le ||g_1||_{L_p(\mathbb{R})} ||w||_{L_1(\mathbb{R})} \le c||w||_{L_1(0,T)}$$

because $||g_1||_{L_n(\mathbb{R})} = ||e^{-k_{\varphi}t}||_{L_n(\mathbb{R}_+)} \le c(k_{\varphi}).$

For h_2 the estimate follows in the same way, setting $g_2(x) = e^{k_{\psi}x}$ to $x \leq 0$ and 0 otherwise. \square In the following, we will denote a generic constant by c and rename it accordingly over the course of a proof. However, the constants in the proofs will never depend on the horizon T. Also, we tacitly use equivalence of norms in \mathbb{R}^2 : $\max\{|a|,|b|\}\approx \sqrt{|a|^2+|b|^2}\approx |a|+|b|$.

LEMMA 4.6. Let (y, λ) solve (4.1) and let Assumption 4.2 hold. Then there exists a constant $c \geq 0$ independent of T, such that

$$(4.8) ||y||_{2\wedge\infty}^2 + ||\lambda||_{2\wedge\infty}^2 \le c \left(||Cy||_{L_2(Y)}^2 + ||R^{-*}B^*\lambda||_{L_2(U)}^2 + ||r||_{1\vee 2}^2 \right),$$

where $||r||_{1\vee 2}^2 := ||l_1||_{1\vee 2}^2 + ||l_2||_{1\vee 2}^2 + ||y_0||^2 + ||\lambda_T||^2$.

Proof. Our first step will be to derive an estimate for ||y(t)|| from (4.6). By Lemma 4.3 we estimate the terms occurring in (4.6) as follows:

(4.9)
$$\int_0^t |\langle K_C C y(s), \varphi(s) \rangle| ds \le ||y(t)|| \frac{M_{\varphi} ||K_C||}{\sqrt{k_{\varphi}}} \sqrt{\int_0^t ||C y(s)||^2 e^{-k_{\varphi}(t-s)} ds},$$

$$(4.10) \quad \int_0^t |\langle R^{-*}B^*\lambda(s), R^{-*}B^*\varphi(s)\rangle| \, ds \leq \|y(t)\| \frac{M_{\varphi}\|BR^{-1}\|}{\sqrt{k_{\varphi}}} \sqrt{\int_0^t \|R^{-*}B^*\lambda(s)\|^2 e^{-k_{\varphi}(t-s)} \, ds},$$

(4.11)
$$\int_0^t \langle l_2(s), \varphi(s) \rangle \, ds \le \|y(t)\| \frac{M_{\varphi}}{\sqrt{k_{\varphi}}} \sqrt{\int_0^t \|l_2(s)\|^2 e^{-k_{\varphi}(t-s)} \, ds},$$

or alternatively:

$$(4.12) \qquad \int_0^t \langle l_2(s), \varphi(s) \rangle \, ds \le \int_0^t \|l_2(s)\| \|\varphi(s)\| \, ds \le \|y(t)\| M_{\varphi} \int_0^t \|l_2(s)\| e^{-k_{\varphi}(t-s)} \, ds,$$

and finally:

$$(4.13) \langle y_0, \varphi(0) \rangle \le ||y_0|| ||y(t)|| M_{\varphi} e^{-k_{\varphi} t} \le ||y_0|| ||y(t)|| M_{\varphi} \sqrt{e^{-k_{\varphi} t}}.$$

Now we substitute all estimates (4.9)-(4.13) into (4.6) while taking the minimum over (4.11) and (4.12) and cancel ||y(t)|| on both sides. Taking squares on both side and using the simple inequality

$$a, b, c, d \ge 0 \Rightarrow (a + b + c + d)^2 \le 4(a^2 + b^2 + c^2 + d^2)$$

we obtain the following pointwise estimate for y:

$$(4.14) ||y(t)||^{2} \leq c \left(\int_{0}^{t} (||Cy(s)||^{2} + ||R^{-*}B\lambda(s)||^{2}) e^{-k_{\varphi}(t-s)} dt + ||y_{0}||^{2} e^{-k_{\varphi}t} \right)$$

$$+ c \min \left\{ \int_{0}^{t} ||l_{2}(s)||^{2} e^{-k_{\varphi}(t-s)} ds, \left(\int_{0}^{t} ||l_{2}(s)|| e^{-k_{\varphi}(t-s)} ds \right)^{2} \right\}.$$

To derive an estimate for $||y||_{C(X)}^2$ we replace all exponential functions by 1, extend the domains of integration from [0,t] to [0,T], and take the maximum over all $t \in [0,T]$:

$$(4.15) ||y||_{C(X)}^2 \le c(||Cy||_{L_2(Y)}^2 + ||R^{-*}B^*\lambda||_{L_2(U)}^2 + \min\{||l_2||_{L_2(X)}^2, ||l_2||_{L_1(X)}^2\} + ||y_0||^2).$$

Similar via (4.7) we obtain:

To obtain an estimate for $||y||_{L_2(X)}^2$ we have to integrate (4.14) over [0, T] and apply Lemma 4.5 to the integral terms in (4.14). Setting:

$$w(s) := ||l_2(s)||^2 \quad \Rightarrow \quad h_1(t) = \int_0^t ||l_2(s)||^2 e^{-k_{\varphi}(t-s)} \, ds,$$

we obtain with Lemma 4.5:

$$\int_0^T \int_0^t \|l_2(s)\|^2 e^{-k_{\varphi}(t-s)} \, ds \, dt = \|h_1\|_{L_1(0,T)} \le c \|w\|_{L_1(0,T)} = c \|l_2\|_{L_2(X)}^2$$

and similarly:

$$\int_0^T \int_0^t (\|Cy(s)\|^2 + \|R^{-*}B^*\lambda(s)\|^2) e^{-k_{\varphi}(t-s)} ds dt \le c(\|Cy(s)\|_{L_2(Y)}^2 + \|R^{-*}B^*\lambda(s)\|_{L_2(U)}^2).$$

If we set

$$w(s) := ||l_2(s)|| \quad \Rightarrow \quad h_1(t) = \int_0^t ||l_2(s)|| e^{-k_{\varphi}(t-s)} \, ds,$$

we obtain

$$\int_0^T \left(\int_0^t \|l_2(s)\| e^{-k_{\varphi}(t-s)} \right)^2 \, ds \, dt = \|h_1\|_{L_2(0,T)}^2 \le c \|w\|_{L_1(0,T)}^2 = c \|l_2\|_{L_1(X)}^2.$$

This yields the desired L_2 -estimate:

$$(4.17) ||y||_{L_2(X)}^2 \le c(||Cy(s)||_{L_2(Y)}^2 + ||R^{-*}B^*\lambda(s)||_{L_2(U)}^2 + \min\{||l_2||_{L_2(X)}^2, ||l_2||_{L_1(X)}^2\} + ||y_0||^2).$$
 In the same way we we obtain

$$(4.18) \quad \|\lambda\|_{L_2(X)}^2 \le c(\|Cy(s)\|_{L_2(Y)}^2 + \|R^{-*}B^*\lambda(s)\|_{L_2(U)}^2 + \min\{\|l_1\|_{L_2(X)}^2, \|l_1\|_{L_1(X)}^2\} + \|\lambda_T\|^2).$$

Now we take the maximum of (4.15) and (4.17) and add it to the maximum of (4.16) and (4.18). Using the definition of the norms

$$\|v\|_{2\wedge\infty}^2 = \max\{\|v\|_{L_2(X)}^2, \|v\|_{C(X)}^2\}, \quad \|w\|_{1\vee 2}^2 = \min\{\|w\|_{L_1(X)}^2, \|w\|_{L_2(X)}^2\}$$

our result follows.

The first two terms on the right-hand side of (4.8) still depend on the state and the adjoint. We therefore present the following representation formula, motivated by [29, Proof of Theorem 1].

LEMMA 4.7. Let (y, λ) solve (4.1). Then

$$(4.19) \quad ||Cy||_{L_2(Y)}^2 + ||R^{-*}B^*\lambda||_{L_2(U)}^2 = \langle \lambda_T, y(T) \rangle - \langle y_0, \lambda(0) \rangle + \int_0^T \langle l_2(s), \lambda(s) \rangle - \langle l_1(s), y(s) \rangle \, ds$$

$$\leq c \left(||\lambda_T|| ||y(T)|| + ||y_0|| ||\lambda(0)|| + ||l_2||_{1\vee 2} ||\lambda||_{2\wedge \infty} + ||l_1||_{1\vee 2} ||y||_{2\wedge \infty} \right).$$

Proof. We apply Lemma 3.4 to the state and adjoint equation, which yields

$$\langle \lambda_T, y(T) \rangle - \langle y_0, \lambda(0) \rangle = \int_0^T \langle l_2(s), \lambda(s) \rangle - \langle l_1(s), y(s) \rangle - \|Cy(s)\|^2 - \|R^{-*}B^*\lambda(s)\|^2 ds.$$

Rearranging the terms yields the result.

THEOREM 4.8. Let Assumption 4.2 hold. Then there is $c \geq 0$ independent of T such that

$$||M^{-1}||_{((L_1(X),||\cdot||_{1\vee 2})\times X)^2\to (C(X),||\cdot||_{2\wedge\infty})^2}\leq c.$$

Proof. Consider $z := (y, \lambda) \in C(X)^2$ and $r := (l_1, \lambda_T, l_2, y_0) \in (L_1(X) \times X)^2$ which satisfy (4.1). Thus, as shown in Lemma 4.6 the estimate (4.8) applies. We substitute (4.19) into (4.8) and apply the Hölder-like inequality of Lemma 4.1 to the integral terms:

$$\begin{aligned} \|z\|_{2\wedge\infty}^2 &:= \|y\|_{2\wedge\infty}^2 + \|\lambda\|_{2\wedge\infty}^2 \\ &\leq c \left(\|\lambda_T\| \|y(T)\| + \|y_0\| \|\lambda(0)\| + \|l_2\|_{1\vee2} \|\lambda\|_{2\wedge\infty} + \|l_1\|_{1\vee2} \|y\|_{2\wedge\infty} + \|r\|_{1\vee2}^2 \right) \\ &\leq c \left((\|\lambda_T\| + \|l_1\|_{1\vee2}) \|y\|_{2\wedge\infty} + (\|y_0\| + \|l_2\|_{1\vee2}) \|\lambda\|_{2\wedge\infty} + \|r\|_{1\vee2}^2 \right) \\ &\leq c \left(\|r\|_{1\vee2} \|z\|_{2\wedge\infty} + \|r\|_{1\vee2}^2 \right). \end{aligned}$$

Application of the simple estimate $c||r||_{1\vee 2}||z||_{2\wedge\infty} \leq \frac{1}{2}(c^2||r||_{1\vee 2}^2 + ||z||_{2\vee\infty}^2)$ yields

$$||M^{-1}r||_{2\wedge\infty} = ||z||_{2\wedge\infty} \le c||r||_{1\vee 2}$$

and thus the desired result.

5. Extension to boundary control with admissible control operators. In this section, we extend the results of Section 3 and 4 to the case of a control operator B which is unbounded as mapping into X but admissible for the semigroup generated by Λ in the sense of [31, Chapter 4]. We refer to [31, Chapter 4] for an in-depth introduction to this topic. The case of an unbounded operator often arises in the case of boundary control. As a consequence of this the operator norm $\|B\|_{U\to X} = \|B^*\|_{X\to U}$ is no longer finite. Additionally, we will allow for a state feedback operator K_B in Assumption 4.2 which is unbounded but admissible. Inspection of the proofs in Section 3 and 4 yields that these operator norms are only used in inequality (4.10) and the dual version of (4.9). In addition, the estimate of the control in Theorem 3.2 is performed via $R^*Ru = B^*\lambda$ using $\|B^*\|_{X\to U}$ in (3.9). All remaining estimates and constants do not involve norms of B or K_B .

The goal of this section is to replace boundedness of B and the feedback operator K_B in Assumption 4.2 by a weaker property, which is known as admissibility. Our strategy of proof will be to show surrogates for (4.10) and the dual version of (4.9) involving K_B , which allow us to generalize our main results to the case of admissible control and feedback operators. The remaining steps of the proofs remain unchanged. The sensitivity results in Theorem 3.2 for the state and adjoint directly carry over as stated in Theorem 5.8, whereas the estimate for the control in the proof of Theorem 3.2 involves the norm of the control operator. We therefore present a modification of this in the proof of Theorem 5.8. However, one only obtains an integral estimate but no uniform estimate on the control in Theorem 3.2. This is the only price to pay for going from bounded to unbounded but admissible control operators.

5.1. Admissible control and observation operators. For ease of reading we recall the definition of admissible control and observation operators. Let $\Lambda: D(\Lambda) \subset X \to X$ be the generator of a semigroup $(\mathcal{T}(t))_{t\geq 0}$ on X. Moreover, let Λ^* be the adjoint operator of Λ with domain $D(\Lambda^*)$. Let X_1 be $D(\Lambda)$ equipped with the norm $\|\cdot\|_1 = \|(\beta I - \Lambda) \cdot \|$ for $\beta \in \rho(\Lambda)$, where $\rho(\Lambda) := \{\beta \in \mathbb{C} \mid \beta I - \Lambda \text{ is invertible and } (\beta I - \Lambda)^{-1} \in L(X, X)\}$ is the resolvent set of Λ . Furthermore, again for $\beta \in \rho(\Lambda)$, we define X_{-1} to be the completion of X with respect to the norm $\|\cdot\|_{-1} = \|(\beta I - \Lambda)^{-1} \cdot \|$. We note that the norms $\|\cdot\|_1$ for different β are equivalent, see [31, Proposition 2.10.1], and the same also holds true for $\|\cdot\|_{-1}$, see [31, Proposition 2.10.2]. Furthermore, by e.g. [31, Proposition 2.10.4], the semigroup $(\mathcal{T}(t))_{t\geq 0}$ can be extended to a semigroup on X_{-1} .

Definition 5.1. ([31, Definition 4.2.1, Definition 4.3.1])

i) $B \in L(U, X_{-1})$ is called an admissible control operator for the semigroup $(\mathcal{T}(t))_{t\geq 0}$, if for some $\tau > 0$, Ran $\Phi_{\tau} \subset X$, where for $u \in L_2(0, \infty; U)$,

$$\Phi_{\tau}u := \int_0^{\tau} \mathcal{T}(\tau - s)Bu(s) \, ds.$$

ii) $C \in L(X_1, Y)$ is called an admissible observation operator for the semigroup $(\mathcal{T}(t))_{t\geq 0}$ if for some $\tau > 0$, Ψ_{τ} has a continuous extension to X, where for $z_0 \in X_1$,

$$(\Psi_{\tau}z_0)(t) := \begin{cases} C\mathcal{T}(t)z_0 & \text{for } t \in [0,\tau] \\ 0 & \text{for } t > \tau. \end{cases}$$

Note, that if i) and ii) in Definition 5.1 are satisfied for one $\tau \geq 0$, they hold for all $\tau \geq 0$, see [31, Proposition 4.2.2, Proposition 4.3.2]. We briefly recall properties of admissible control operators, which are important in the remainder of this section.

PROPOSITION 5.2. Let B be an admissible control operator for the semigroup $(\mathcal{T}(t))_{t>0}$. Then,

- i) B^* is an admissible observation operator for the adjoint semigroup $(\mathcal{T}^*(t))_{t>0}$.
- ii) For all $t \geq 0$ and $y(t) \in X$ there exists a constant $K_t \geq 0$, such that

(5.1)
$$\int_0^t \|B^* \mathcal{T}^*(t-s) y(t)\|^2 ds \le K_t^2 \|y(t)\|^2.$$

iii) If $(\mathcal{T}^*(t))_{t>0}$ is exponentially stable, the constant K_t can be chosen independently of t.

Proof. i) follows by the duality result [31, Theorem 4.4.3]. For ii), see [31, Definition 4.3.1]. The fact that the bound can be chosen independently of t, as stated in iii), follows by [31, Remark 4.3.5].

5.2. Well-posed linear systems. In addition to the concept of admissible control and observation operators as introduced in the previous chapter, we will make use of the notion of well-posed linear systems, for which certain desirable properties hold. In particular these systems enjoy continuity of the state trajectory and boundedness of input-to-state, state-to-output and input-to-output maps as partly defined in Definition 5.1, although the generating operators can be unbounded. For an in-depth treatment of this topic, we refer the reader to the seminal papers [26, 27, 33] to the monograph [28] and the survey articles [34, 32]. A possible approach of defining well-posed linear systems is to require particular properties of the above mentioned maps, cf. [32, Definition 3.1]. Properties include, e.g., that the initial value to state map is given by a strongly continuous semigroup, that the state does not depend on the future input and that the past output does not depend on the future input, the latter two often called causality. It turns out, that the operators (Λ, B, C) generating such an well-posed linear systems are precisely characterized by Λ generating a semigroup, admissibility of B and C with respect to the semigroup, and a condition on the transfer function, assuring boundedness of the input-output map, cf. [32, Proposition 4.9] and [28, Theorem 4.2.1, Theorem 4.4.2].

PROPOSITION 5.3. ([32, Proposition 4.9]) A triple of operators (A, B, C) is well-posed on (U, X, Y) iff the following conditions hold:

- A generates a C0-semigroup $(\mathcal{T}(t))_{t>0}$ on X,
- $\mathcal{B} \in L(U, X_{-1})$ is an admissible control operator for $(\mathcal{T}(t))_{t>0}$,
- $C \in L(X_1, Y)$ is an admissible observation operator for $(T(t))_{t\geq 0}$,
- some (hence every) transfer function associated with (A, B, C) is proper,

where the transfer function is characterized by $G(s) - G(\beta) = \mathcal{C}((sI - \mathcal{A})^{-1} - (\beta I - \mathcal{A})^{-1})\mathcal{B}$ for every $s, \beta \in \{s \in \mathbb{C} \mid \Re(s) > \omega\}$, where ω is the growth bound of the semigroup. An analytic function with a domain in some right-half plane is called proper if it is bounded on some right-half plane [32, p.8].

The optimality conditions (2.3) are derived under the assumption of boundedness of the control and observation operator via [20, Chapter 4]. To the authors best knowledge, optimality conditions in this abstract setting with unbounded control and observation are not avaliable. Thus, in the remainder, we will assume that (y, u, λ) solve the extremal equations, rather than (y, u) solves the optimal control problem. A unique solution of the extremal equations can be obtained in the framework of well posed systems, as state in the following theorem:

THEOREM 5.4. Let (Λ, B, C) form a well-posed system. Then (2.3) has a unique solution $(y, \lambda) \in C(X)^2$ which satisfy the equations in the mild sense.

Proof. It suffices to show that $\begin{pmatrix} C^*C & \Lambda^* \\ \Lambda & BQ^{-*}B^* \end{pmatrix}$ generates a C0-semigroup. Defining

$$\mathcal{A} = \begin{pmatrix} 0 & \Lambda^* \\ \Lambda & 0 \end{pmatrix}, \, \mathcal{B} = \begin{pmatrix} 0 & C^* \\ B & 0 \end{pmatrix}, \, \mathcal{C} = \begin{pmatrix} 0 & B^* \\ C & 0 \end{pmatrix}$$

by the the well posedness of (Λ, B, C) it follows that $(\mathcal{A}, \mathcal{B}, \mathcal{C})$ forms a well-posed system. Choosing the bounded feedback operator $K = \begin{pmatrix} I & 0 \\ 0 & Q^{-1} \end{pmatrix}$, the closed loop system $(\mathcal{A} + \mathcal{B}K\mathcal{C}, \mathcal{B}, \mathcal{C})$ with

$$\mathcal{A} + \mathcal{B}K\mathcal{C} = \begin{pmatrix} C^*C & \Lambda^* \\ \Lambda & BQ^{-*}B^* \end{pmatrix} \text{ forms another well-posed system, cf. [28, Theorem 7.1.2] or [33].}$$

In particular, $\mathcal{A} + \mathcal{B}K\mathcal{C}$ generates a C0-semigroup and $\mathcal{C} = \begin{pmatrix} 0 & B^* \\ C & 0 \end{pmatrix}$ is an admissible observation operator for the semigroup generated by $\mathcal{A} + \mathcal{B}K\mathcal{C}$.

Assumption 5.5. Let Λ generate a C0-semigroup $(\mathcal{T}(t)_{t\geq 0} \text{ on } X, C \in L(X,Y) \text{ and } B \in L(U,X_{-1})$ be an admissible control operator for $(\mathcal{T}(t))_{t\geq 0}$. Further assume that

- 1. (Λ, C) is exponentially detectable, i.e., there exists $K_C \in L(Y, X)$ such that the semigroup generated by $\Lambda^* + C^*K_C^*$ is exponentially stable.
- 2. (Λ, B) is exponentially stabilizable, i.e., there exists $K_B \in L(X_1, U)$
 - (Λ, B, K_B) well-posed on (U, X, U), and
 - the semigroup generated by $\Lambda + BK_B$ is exponentially stable.

Remark 5.6. In a similar fashion, one could allow for unbounded but admissible C and K_C when assuming B and K_B to be bounded. The case where all of the operators C, B, K_C and K_B are unbounded but admissible cannot be included in all generality. This is due to the non-existence of perturbation results for this case, i.e., not every admissible observation operator C for Λ is admissible for $\Lambda + BK$, if B is an admissible control operator and K is an unbounded but admissible state feedback operator. Mixed perturbation results of this kind are a very delicate matter and we refer to [31, Proposition 5.5.2, Example 5.5.3, Proposition 10.1.10] and the discussion after [31, Corollary 5.5.1] for results on that subject. Moreover, we refer to the example in [31, Section 10.8] where such a perturbation result holds. Finally, for the particular case of K = LC with $L \in L(Y, U)$, such perturbation results are established in [28, Theorem 7.1.2] or [33].

The strategy now consists of applying the admissibility estimate of Proposition 5.2 to φ and ψ in Lemma 4.3. As the underlying semigroup is exponentially stable, we may apply (5.1) to the respective admissible control operators with K_t independent of t. This will turn out as an appropriate replacement for the assumption $||B^*||, ||K_B|| < \infty$.

5.3. Bounds for the solution operator. We will now use these properties to derive the desired bound on the operator norm in the case of unbounded but admissible control operators.

Theorem 5.7. Consider M defined in (2.3) and let Assumption 5.5 hold. Then there is $c \ge 0$ independent of T, such that

$$\|M^{-1}\|_{((L_1(X),\|\cdot\|_{1\vee 2})\times X)^2\to (C(X),\|\cdot\|_{2\wedge\infty})^2}\leq c.$$

Proof. As already noted, the only step in the proof of Theorem 4.8 where the operator norm of B and K_B is needed, is the proof of (4.10) and (4.9) with their respective dual counterparts. It is thus sufficient to show a modification of this inequality which circumvents this operator norm

estimate by using exponential stability of the test functions and the fact that B is an admissible control operators. We will sketch the proof only for the case of B as the estimate for K_B in the dual counterpart of (4.9) follows completely analogously.

Let φ solve (4.2), i.e., $\varphi(s) = \mathcal{T}_{cl}^*(t-s)y(t)$ where, as above, $(\mathcal{T}_{cl}^*(t))_{t\geq 0}$ is the semigroup generated by $(\Lambda^* + C^*K_C^*)$. As by our assumption B is an admissible control operator for the semigroup generated by Λ , and $K_CC \in L(X,X)$, we obtain by a perturbation result, cf. [31, Theorem 5.4.2], that B is an admissible control operator for $(\Lambda + K_CC)$. Hence B^* is an admissible observation operator for the adjoint semigroup generated by $(\Lambda + K_CC)^*$. We show that the critical estimate (4.10) still holds with different constants which do not involve the operator norm of B. First, a simple calculation using the fact that R^{-*} is bounded from U to U and the Cauchy-Schwarz inequality twice yields

$$(5.2) \int_{0}^{t} |\langle R^{-*}B^{*}\lambda(s), R^{-*}B^{*}\varphi(s)\rangle| \, ds \leq c \int_{0}^{t} e^{-\frac{k\varphi}{2}(t-s)} \|R^{-*}B^{*}\lambda(s)\| \|B^{*}e^{\frac{k\varphi}{2}(t-s)}\varphi(s)\| \, ds$$

$$\leq c \sqrt{\int_{0}^{t} e^{-k\varphi(t-s)} \|R^{-*}B^{*}\lambda(s)\|^{2} \, ds} \sqrt{\int_{0}^{t} \|B^{*}e^{\frac{k\varphi}{2}(t-s)}\varphi(s)\|^{2} \, ds}.$$

Using the exponential stability $\|\mathcal{T}_{\mathrm{cl}}^*(t-s)\| \leq M_{\varphi}e^{-k_{\varphi}(t-s)}$, the scaled semigroup $(e^{\frac{k_{\varphi}}{2}(t-s)}\mathcal{T}_{\mathrm{cl}}^*(t))_{t\geq 0}$ is still exponentially stable, and hence we employ Proposition 5.2 ii) and iii) and get

$$\sqrt{\int_0^t \|B^* e^{\frac{k_{\varphi}}{2}(t-s)} \varphi(s)\|^2 \, ds} = \sqrt{\int_0^t \|B^* e^{\frac{k_{\varphi}}{2}(t-s)} \mathcal{T}_{\mathrm{cl}}^*(t-s) y(t)\|^2 \, ds} \leq K \|y(t)\|$$

with K independent of t. Inserting this into (5.2), we conclude

$$\int_0^t |\langle R^{-*}B^*\lambda(s), R^{-*}B^*\varphi(s)\rangle| \, ds \le c \|y(t)\| \sqrt{\int_0^t e^{-k_\varphi(t-s)} \|R^{-*}B^*\lambda(s)\|^2 \, ds},$$

which yields the desired replacement for (4.10) with a different constant independent of the norm of B. As the remaining results in Section 4, namely Lemma 4.6, Lemma 4.7 and Theorem 4.8 do not use the boundedness, we can conclude the result analogously to the bounded case.

As a consequence of Theorem 5.7, the estimates in Theorem 3.2 also hold true in the case of unbounded control and feedback operators with constants independent of the horizon T, except for the uniform estimate for the control. This is the statement of the following theorem.

THEOREM 5.8. Let Assumption 5.5 hold. Assume $(\delta y, \delta \lambda) \in C(X)^2$ solves (3.1) with $\varepsilon_1, \varepsilon_2 \in L_1(X)$. Then there is a constant $c \geq 0$ and a scaling factor $\mu > 0$, both independent of T, such that defining

$$\rho := \|e^{-\mu t}\varepsilon_1(t)\|_E + \|e^{-\mu t}\varepsilon_2(t)\|_E$$

for $E := L_1(X)$ or $E := L_2(X)$, we have

$$||e^{-\mu t}\delta y||_{L_2(X)} + ||e^{-\mu t}\delta u||_{L_2(U)} + ||e^{-\mu t}\delta \lambda||_{L_2(X)} \le c\rho,$$
$$||e^{-\mu t}\delta y||_{C(X)} + ||e^{-\mu t}\delta \lambda||_{C(X)} \le c\rho.$$

Proof. First, choosing μ , such that $\mu \| M^{-1} \|_{(L_2(X) \times X)^2 \to L_2(X)^2 < 1}$, we conclude by the same reasoning as in the proof of Theorem 3.2 the estimates for the state and adjoint

with $c \geq 0$ independent of T as the occurring operator norms can be bounded by a constant $c \geq 0$ independently of T by Theorem 5.7. To estimate the control, we compute, setting $\widetilde{\delta y}(t) := e^{-\mu t} \delta y(t)$, $\widetilde{\delta \lambda}(t) := e^{-\mu t} \delta \lambda(t)$ and $\widetilde{\delta u}(t) := e^{-\mu t} \delta u$, that

$$\|\widetilde{\delta u}\|_{L_2(U)}^2 = \int_0^T \|Q^{-1}B^*\widetilde{\delta \lambda}(t)\|^2 dt \le c \int_0^T \|C\widetilde{\delta y}(t)\|^2 + \|R^{-*}B^*\widetilde{\delta \lambda}(t)\|^2 dt.$$

Similar to Lemma 4.7, now for the scaled system (3.2), using Lemma 3.4 for the state and adjoint equation, we obtain with $\tilde{\varepsilon}_i(t) = e^{-\mu t} \varepsilon_i(t)$, i = 1, 2 using the Hölder inequality of Lemma 4.1

$$\int_{0}^{T} \|C\widetilde{\delta y}(t)\|^{2} + \|R^{-*}B^{*}\widetilde{\delta\lambda}(t)\|^{2} dt = \int_{0}^{T} -\langle \widetilde{\varepsilon}_{1}(t), \widetilde{\delta y}(t) \rangle + \langle \widetilde{\varepsilon}_{2}(t), \widetilde{\delta\lambda}(t) \rangle - 2\mu \langle \widetilde{\delta y}(t), \widetilde{\delta\lambda}(t) \rangle dt$$

$$\leq (\|\widetilde{\varepsilon}_{1}\|_{1\vee 2} + \|\widetilde{\varepsilon}_{2}\|_{1\vee 2})(\|\widetilde{\delta y}\|_{2\wedge \infty} + \|\widetilde{\delta\lambda}\|_{2\wedge \infty}) + 2\mu \|\widetilde{\delta y}\|_{L_{2}(X)} \|\widetilde{\delta\lambda}\|_{L_{2}(X)}$$

$$\leq c\rho^{2}$$

where in the last estimate we used the assumption on $\tilde{\varepsilon_i}$, i = 1, 2 and the estimate on the state and the adjoint (5.3). Taking the square root yields the result for the control.

Remark 5.9. If we have additional structure, the results of Theorem 4.8 and Theorem 5.7 can be refined by bootstrapping arguments. For instance, in the parabolic case there is an additional, more regular space with the continuous dense embedding $V \hookrightarrow X$, such that Λ can be extended to an operator $\bar{\Lambda} \in L(V, V^*)$ and we have a Gårding inequality for $y \in L_2(0, T; V)$:

$$(5.4) \exists \omega \in \mathbb{R}, \alpha > 0: \quad \alpha \|y\|_{L_2(V)}^2 \le \int_0^T -\langle \bar{\Lambda} y(t), y(t) \rangle_{V^* \times V} \, dt + \omega \|y\|_{L_2(X)}^2.$$

In this particular case, the state equation yields the regularity $y \in L_2(0,T;V)$, $y' \in L_2(0,T;V^*)$ and we compute, using integration by parts on $\langle y',y\rangle_{V^*\times V}$, (4.19), and assuming boundedness of $||B^*||_{V\to U}$ (which is weaker than boundedness of $||B^*||_{X\to U}$):

$$\int_{0}^{T} -\langle \bar{\Lambda}y(t), y(t) \rangle_{V^{*} \times V} dt = \int_{0}^{T} -\langle y'(t), y(t) \rangle_{V^{*} \times V} + \langle l_{2}(t) + B(R^{*}R)^{-1}B^{*}\lambda(t), y(t) \rangle_{X} dt$$

$$\leq \frac{1}{2} \left(\|y(0)\|^{2} - \|y(T)\|^{2} \right) + \|l_{2}\|_{1 \vee 2} \|y\|_{2 \wedge \infty} + \|R^{-*}B^{*}\lambda\|_{L_{2}(U)} \|R^{-*}B^{*}y\|_{L_{2}(U)}$$

$$\leq \frac{1}{2} \|y_{0}\|^{2} + \|r\|_{1 \vee 2} \|z\|_{2 \wedge \infty} + c(\|r\|_{1 \vee 2} + \|z\|_{2 \wedge \infty}) \|R^{-*}B^{*}\|_{V \to U} \|y\|_{L_{2}(V)}.$$

By Theorem 5.8 we may estimate $||z||_{2\wedge\infty}$ by $||r||_{1\vee 2}$ and substitute the result into (5.4). A short computation yields $||y||_{L_2(V)} \leq c||r||_{1\vee 2}$ and similarly $||\lambda||_{L_2(V)} \leq c||r||_{1\vee 2}$. Hence, there is a T-independent bound:

$$||M^{-1}||_{(L_1(X),||\cdot||_{1\vee 2}\times X)^2\to L_2(V)^2} \le c.$$

By further bootstrapping via:

$$y' = \bar{\Lambda}y + B(R^*R)^{-1}B^*\lambda + l_2$$
 in $L_2(V^*)$

we obtain $||y'||_{L_2(V^*)} \le c(||r||_{1\vee 2} + ||l_2||_{L_2(X)})$ and similarly an estimate for $||\lambda'||_{L_2(V^*)}$. Thus, also for the parabolic space W([0,T]), which is equipped with the norm $||v||_{W([0,T])} = ||v||_{L_2(V)} + ||v'||_{L_2(V^*)}$, we get the T-independent bound:

$$||M^{-1}||_{(L_2(X)\times X)^2\to W([0,T])^2}\le c.$$

These additional estimates can be used to obtain results in Theorem 3.2 and Theorem 5.8, and also in Theorem 6.2 below in the respective additional norms.

6. An exponential turnpike result. In this section, we extend the scaling approach employed in Theorem 3.2 to deduce a turnpike result. It is an adaption of existing results for parabolic equations in [15, Theorem 5.2, Corollary 5.3] and [5, Lemma 15]. In the case of optimal control problems governed by general evolution equations in Hilbert spaces, turnpike theorems were given in [29] for C_0 -semigroups with bounded control and observation operators, and for boundary controlled parabolic equations. Here we extend these results to the boundary control of non-parabolic equations with admissible control operators. First, we introduce the solution of a steady state optimization problem, namely $(\bar{y}, \bar{\lambda}) \in X^2$ being a minimizer of

(6.1)
$$\min_{\bar{y},\bar{u}} \frac{1}{2} \|C(\bar{y} - y_d)\|_Y^2 + \frac{1}{2} \|R(\bar{u} - u_d)\|_U^2$$
s.t.
$$-\Lambda \bar{y} - B\bar{u} = 0,$$

or, equivalently, as the problem is convex, by coercivity of R, $(\bar{y}, \bar{\lambda})$ solve the corresponding first order conditions

(6.2)
$$\begin{pmatrix} C^*C & -\Lambda^* \\ -\Lambda & -BQ^{-1}B^* \end{pmatrix} \begin{pmatrix} \bar{y} \\ \bar{\lambda} \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ Bu_d \end{pmatrix},$$

where $Q = R^*R$ and we eliminated the control via $\bar{u} = Q^{-1}B^*\bar{\lambda} + u_d$.

LEMMA 6.1. Let (y, u, λ) solve (2.2) with f = 0. Moreover, let $(\bar{y}, \bar{u}, \bar{\lambda})$ solve the corresponding steady state problem (6.2). Then $(\delta y, \delta \lambda) := (y - \bar{y}, \lambda - \bar{\lambda})$ solves

(6.3)
$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} \delta y \\ \delta \lambda \end{pmatrix} = \begin{pmatrix} 0 \\ -\bar{\lambda} \\ 0 \\ y_0 - \bar{y} \end{pmatrix},$$

where $E_0\delta y := \delta y(0)$ and $E_T\delta \lambda := \delta \lambda(T)$ and $\delta u := u - \bar{u} = Q^{-1}B^*\delta \lambda$.

Proof. Using (6.2) and $\frac{d}{dt}\bar{y} = \frac{d}{dt}\bar{\lambda} = 0$ yields that $(\bar{y}, \bar{\lambda})$ satisfies

(6.4)
$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} \bar{y} \\ \bar{\lambda} \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ \bar{\lambda} \\ Bu_d \\ \bar{y} \end{pmatrix}$$

Considering the linear system for (y, u, p) solving (2.2) with f = 0, i.e.,

(6.5)
$$\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix} \begin{pmatrix} y \\ \lambda \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ 0 \\ Bu_d \\ y_0 \end{pmatrix}$$

we conclude the result by subtracting (6.4) from (6.5).

THEOREM 6.2. Let Assumption 5.5 hold. Assume (y, u, λ) solves (2.2). Moreover, let $(\bar{y}, \bar{u}, \bar{\lambda})$ solve the corresponding steady state problem (6.1). Then, defining $(\delta y, \delta u, \delta \lambda) := (y - \bar{y}, u - \bar{u}, \lambda - \bar{\lambda})$, there exist $\mu > 0$ and a constant $c \ge 0$, both independent of T, such that

(6.6)
$$\|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta y\|_{L_{2}(X)} + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta u\|_{L_{2}(U)} +$$

$$+ \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta \lambda\|_{L_{2}(X)} \le c \left(\|y_{0} - \bar{y}\| + \|\bar{\lambda}\|\right),$$
(6.7)
$$\|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta y\|_{C(X)} + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta \lambda\|_{C(X)} \le c \left(\|y_{0} - \bar{y}\| + \|\bar{\lambda}\|\right).$$

Proof. By Lemma 6.1, $(\delta y, \delta u, \delta \lambda)$ solves (6.3). By Theorem 5.8, $\|M^{-1}\|_{(L_2(X) \times X)^2 \to L_2(X)^2}$ can be bounded independently of T. If we introduce a scaling factor $0 < \mu < \frac{1}{\|M^{-1}\|_{(L_2(X) \times X)^2 \to L_2(X)^2}}$ and scaled variables $\widetilde{\delta y} := \frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta y$, $\widetilde{\delta u} := \frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta u$ and $\widetilde{\delta \lambda} := \frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta \lambda$, we can compute

(6.3)
$$\iff \underbrace{\left(\underbrace{\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ 0 & E_T \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix}}_{=:M} + \mu \underbrace{\begin{pmatrix} 0 & F \\ 0 & 0 \\ -F & 0 \\ 0 & 0 \end{pmatrix}}_{=:\widetilde{P}} \right) \begin{pmatrix} \widetilde{\delta y} \\ \widetilde{\delta \lambda} \end{pmatrix} = \frac{1}{1 + e^{-\mu T}} \begin{pmatrix} 0 \\ -\bar{\lambda} \\ 0 \\ y_0 - \bar{y} \end{pmatrix}$$

where $F:=\frac{e^{-\mu(T-t)}-e^{-\mu t}}{e^{-\mu t}+e^{-\mu(T-t)}}$ and the factor $\frac{1}{1+e^{-\mu T}}$ arises due to the scaling of the initial values. The proof for the estimate of the state and the adjoint in (6.6) and (6.7) is similar to the one of Theorem 3.2. Defining $\tilde{z}:=(\widetilde{\delta y},\widetilde{\delta\lambda})$ and $\tilde{r}:=\frac{1}{1+e^{-\mu T}}(0,-\bar{\lambda},0,y_0-\bar{y})$, we get

$$(6.8) \qquad (M+\mu\widetilde{P})\tilde{z}=\tilde{r} \quad \Rightarrow \quad (I+\mu\widetilde{P}M^{-1})M\tilde{z}=\tilde{r} \quad \Rightarrow \quad \tilde{z}=M^{-1}(I+\mu\widetilde{P}M^{-1})^{-1}\tilde{r}.$$

Again, as in the proof of Theorem 3.2 by a standard Neumann series argument, cf. [19, Theorem 2.14], as we chose μ such that $\beta := \mu \|M^{-1}\|_{(L_2(X)\times X)^2\to L_2(X)^2} < 1$, $(I+\mu \widetilde{P}M^{-1})$ is invertible as $\|F\|_{L_2(X)\to L_2(X)} \le 1$ and thus $\|\widetilde{P}\|_{L_2(X)^2\to (L_2(X)\times X)^2} \le 1$. Therefore, using the Neumann series representation of $(I+\mu \widetilde{P}M^{-1})^{-1}$ yields

Thus, we conclude with (6.8) and (6.9)

(6.10)
$$\|\tilde{z}\|_{L_{2}(X)^{2}} \leq \frac{\|M^{-1}\|_{(L_{2}(X)\times X)^{2}\to L_{2}(X)^{2}}}{1-\beta} \|\tilde{r}\|_{(L_{2}(X)\times X)^{2}},$$

$$\|\tilde{z}\|_{C(X)^{2}} \leq \frac{\|M^{-1}\|_{(L_{2}(X)\times X)^{2}\to C(X)^{2}}}{1-\beta} \|\tilde{r}\|_{(L_{2}(X)\times X)^{2}}.$$

Finally, $\|\tilde{r}\|_{(L_2(X)\times X)^2} \leq \|y_0 - \bar{y}\| + \|\bar{\lambda}\|$ and going back to the original variables yields the result (6.7).

The remainder of this proof will now consist of estimating the control to conclude (6.6). To this end, we recall the approach taken in the proof of Theorem 5.8 and compute

$$(6.11) \|\widetilde{\delta u}\|_{L_2(U)}^2 = \int_0^T \|Q^{-1}B^*\widetilde{\delta \lambda}(t)\|^2 dt \le c \int_0^T \|C\widetilde{\delta y}(t)\|^2 + \|R^{-*}B^*\widetilde{\delta \lambda}(t)\|^2 dt.$$

Similar to Lemma 4.7, using Lemma 3.4 for the scaled state and adjoint equation of system, we obtain

$$\int_{0}^{T} \|C\widetilde{\delta y}(t)\|^{2} + \|B^{*}\widetilde{\delta\lambda}(t)\|^{2} dt$$

$$= \langle \widetilde{\delta y}(T), \widetilde{\delta\lambda}(T) \rangle - \langle \widetilde{\delta y}(0), \widetilde{\delta\lambda}(0) \rangle + \mu \int_{0}^{T} \langle F\widetilde{\delta y}(t), \widetilde{\delta\lambda}(t) \rangle + \langle \widetilde{\delta y}(t), F\widetilde{\delta\lambda}(t) \rangle dt$$

$$\leq \|\widetilde{\delta y}(T)\| \|\widetilde{\delta\lambda}(T)\| + \|\widetilde{\delta y}(0)\| \|\widetilde{\delta\lambda}(0)\| + 2\mu \|\widetilde{\delta y}\|_{L_{2}(X)} \|\widetilde{\delta\lambda}\|_{L_{2}(X)}$$

$$\leq (\|\widetilde{\delta y}(T)\| + \|\widetilde{\delta\lambda}(0)\|)(\|\widetilde{\delta y}(0)\| + \|\widetilde{\delta\lambda}(T)\|) + \mu(\|\widetilde{\delta y}\|_{L_{2}(X)}^{2} + \|\widetilde{\delta\lambda}\|_{L_{2}(X)}^{2}),$$
(6.12)

where we used $||F||_{L_2(X)\to L_2(X)} \le 1$. In order to estimate the end time value of the state and the initial value of the adjoint we compute:

$$\|\widetilde{\delta y}(T)\| + \|\widetilde{\delta \lambda}(0)\| = \frac{1}{1 + e^{-\mu T}} (\|\delta y(T)\| + \|\delta \lambda(0)\|) \le \|\delta y\|_{C(X)} + \|\delta \lambda\|_{C(X)}$$
$$\le \|M^{-1}\|_{(L_2(X) \times X)^2 \to C(X)^2} (\|y_0 - \bar{y}\| + \|\bar{\lambda}\|).$$

Inserting this into (6.12) and (6.11), and using the estimate for the state and adjoint (6.10), we get

$$\|\widetilde{\delta u}\|_{L_{2}(U)}^{2} \leq c \left(\|M^{-1}\|_{(L_{2}(X)\times X)^{2}\to C(X)^{2}} \left(\|y_{0} - \bar{y}\|^{2} + \|\bar{\lambda}\|^{2} \right) + \mu \frac{\|M^{-1}\|_{(L_{2}(X)\times X)^{2}\to L_{2}(X)^{2}}^{2}}{(1-\beta)^{2}} \left(\|y_{0} - \bar{y}\|^{2} + \|\bar{\lambda}\|^{2} \right) \right),$$

where taking the square root and using $\mu \| M^{-1} \|_{(L_2(X) \times X)^2 \to L_2(X)^2} < 1$ yields

$$\|\widetilde{\delta u}\|_{L_2(U)} \le c \left(\sqrt{\|M^{-1}\|_{(L_2(X)\times X)^2 \to C(X)^2}} + \frac{\sqrt{\|M^{-1}\|_{(L_2(X)\times X)^2 \to L_2(X)^2}}}{(1-\beta)^2} \right) (\|y_0 - \bar{y}\| + \|\bar{\lambda}\|).$$

Theorem 5.7 yields bounds for the operator norms independently of T. Together with (6.10), this gives the estimate (6.6).

Remark 6.3. Similarly to Theorem 3.2, one could deduce a bound on $\|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}}(u(t) - \bar{u})\|_{L_{\infty}(U)}$, if the control operator is bounded.

We give a short interpretation of the two estimates given in Theorem 6.2. For the first inequality, i.e., (6.6) consider a fixed $0 < \varepsilon < 1/2$ and $t \in [\varepsilon T, (1-\varepsilon)T]$. Then, if $T \to \infty$, the two scaling terms $e^{-\mu t}$ and $e^{-\mu(T-t)}$ approach zero exponentially fast and we estimate

$$\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \geq \frac{1}{e^{-\mu\varepsilon T} + e^{-\mu(T-(1-\varepsilon)T)}} = \frac{2}{e^{-\mu\varepsilon T}}$$

and hence e.g. for the difference of state y and its turnpike \bar{y} ,

$$\int_{0}^{T} \|\frac{y(t) - \bar{y}}{e^{-\mu t} + e^{-\mu(T-t)}}\|^{2} dt \ge \int_{\varepsilon T}^{(1-\varepsilon)T} \|\frac{y(t) - \bar{y}}{e^{-\mu t} + e^{-\mu(T-t)}}\|^{2} dt \ge \frac{4}{e^{-2\mu\varepsilon T}} \int_{\varepsilon T}^{(1-\varepsilon)T} \|y(t) - \bar{y}\|^{2} dt$$

which, using (6.6) of Theorem 6.2, implies

$$\int_{\varepsilon T}^{(1-\varepsilon)T} \|y(t) - \bar{y}\|^2 dt \le c e^{-2\mu\varepsilon T} \left(\|y_0 - \bar{y}\| + \|\bar{\lambda}\|\right)^2.$$

Proceeding analogous for the adjoint and the control, we have

$$\|y - \bar{y}\|_{L_2(\varepsilon T, (1-\varepsilon)T; X)}, \|u - \bar{u}\|_{L_2(\varepsilon T, (1-\varepsilon)T; U)}, \|\lambda - \bar{\lambda}\|_{L_2(\varepsilon T, (1-\varepsilon)T; X)} \to 0, \quad \text{if } T \to \infty$$

i.e., L_2 -convergence on a part of the whole time interval [0,T]. The convergence rate is exponential and the part is growing linearly in T, as its length is $(1-2\varepsilon)T$. Moreover, its fraction of the whole interval is constant, as $(1-2\varepsilon)T/T=1-2\varepsilon$. Thus, on a fixed percentage of the interval [0,T], the state, control and adjoint converge to the turnpike in the L_2 norm, as the horizon T goes to infinity.

For the second inequality, i.e., (6.7), rewriting the pointwise estimate, we have

$$||y(t) - \bar{y}|| + ||\lambda(t) - \bar{\lambda}|| \le c(e^{-\mu t} + e^{-\mu(T-t)})$$

for every $t \in [0, T]$. Therefore, if we fix $0 < \varepsilon < \frac{1}{2}$ and take the maximum over all $t \in [\varepsilon T, (1 - \varepsilon)T]$, we get

$$||y - \bar{y}||_{C(\varepsilon T, (1-\varepsilon)T; X)} + ||\lambda - \bar{\lambda}||_{C(\varepsilon T, (1-\varepsilon)T; X)} \le 2ce^{-\mu\varepsilon T}.$$

The right-hand side approaches zero exponentially fast as $T \to \infty$. Therefore, for each $\varepsilon \in]0, \frac{1}{2}[$, we obtain uniform exponential convergence on the interval $[\varepsilon T, (1-\varepsilon)T]$ of y and λ to \bar{y} and $\bar{\lambda}$ respectively as $T \to \infty$. Again, similar to i), with $(1-2\varepsilon)T/T = (1-2\varepsilon)$, we conclude that on a fixed fraction of the whole interval [0,T], the state and adjoint converge to the turnpike in the maximum norm.

7. Extension to a condition on the terminal state. In this section, we discuss an extension of our result to optimal control problems with a condition on the terminal state, i.e. adding in Problem 2.1 a terminal condition $y(T) = y_T \in X$. In the finite dimensional case, linear and nonlinear initial and terminal conditions were included [30]. In the Hilbert space setting, this problem was discussed in [29, Section 2.6], where however the lack of invertibility of the Lyapunov operator prohibited to derive a result with a terminal state condition. To assure existence of an

optimal solution for arbitrary initial and terminal data, we assume (Λ, B) to be exactly controllable in time t_c where $0 < t_c \le T$, i.e. for any initial datum $y_0 \in X$ and terminal state $y_T \in X$, we can find a control that drives the state from y_0 to y_T in any time $T \ge t_c$. This assumption excludes parabolic equations with control not on the whole domain, but is however fulfilled by many hyperbolic systems, cf. Section 8. For a discussion of controllability topics of PDEs we refer the reader to the overview article [37]. Moreover, a crucial point in deriving Pontryagin Maximum Principles for problems including both initial and terminal conditions on the state in infinite dimensions is a codimensionality condition of the reachable set in X, cf. [20, Chapter 4]. This assumption is automatically satisfied if one assumes exact controllability, as the reachable set is the whole space X. For bounded control and observation operators, the optimal solutions satisfy the dynamics (2.2) with $y(0) = y_0$ and $y(T) = y_T$, where no terminal condition on the adjoint is imposed [20, p.135, Theorem 1.6]. By eliminating the control, we obtain the extremal equations

(7.1)
$$\underbrace{\begin{pmatrix} C^*C & -\frac{d}{dt} - \Lambda^* \\ E_T & 0 \\ \frac{d}{dt} - \Lambda & -BQ^{-1}B^* \\ E_0 & 0 \end{pmatrix}}_{\mathcal{C}} \begin{pmatrix} y \\ \lambda \end{pmatrix} = \begin{pmatrix} C^*Cy_d \\ y_T \\ Bu_d + f \\ y_0 \end{pmatrix}.$$

Now, similar to Theorem 4.8, we obtain a T-independent bound under an exact controllability assumption:

THEOREM 7.1. Let (Λ, B) is exactly controllable in time $0 < t_c \le T$ and (Λ, C) be exponentially detectable. Then

$$\|\widetilde{M}^{-1}\|_{((L_1(X),\|\cdot\|_{1\vee 2})\times X)^2\to (C(X),\|\cdot\|_{2\wedge\infty})^2}\leq c.$$

where c > 0 is independent of T.

Proof. We proceed analogous to Section 4. In order to estimate the unknown $\lambda(T)$ which appears on the right-hand side of (4.16) and (4.18), we use a well known exact observability estimate on the solution λ of the dual system which is equivalent to exact controllability of (Λ, B) , cf. [6, Theorem 2.42] which was proven first in the seminal paper [21] with the Hilbert Uniqueness Method. Considering the adjoint equation

$$-\lambda' - \Lambda^*\lambda + C^*Cy = \varepsilon_1$$

and a straightforward generalization of the method for finite dimensional systems of [23, Remark 2.1] to this infinite dimensional case by exact controllability, we obtain

$$\|\lambda(T)\|^{2} \leq c \left(\int_{T-t_{c}}^{T} \|Cy(t)\|^{2} + \|B^{*}\lambda(t)\|^{2} dt + \|l_{1}\|_{1 \vee 2} \right)$$

with a constant $c(t_c) \geq 0$ independent of T. Thus, the statement (4.8) of Lemma 4.6 holds with an upper bound including only $(\|l_1\|_{1\vee 2}^2 + \|l_2\|_{1\vee 2}^2 + \|y_0\|^2)$. Using Lemma 4.7 and the fact that $y(T) = y_T$ is a datum, a straightforward adaption of the proof of Theorem 4.8, where y_T plays the role of λ_T , yields the result.

The result about limited propagation of perturbations from Theorem 3.2 directly carries over to the case of terminal conditions, whereas the turnpike result Theorem 6.2 reads:

THEOREM 7.2. Let $(y, \lambda) \in C(X)^2$ solve (7.1) and $(\bar{y}, \bar{\lambda}) \in X^2$ solve the corresponding steady state problem and set $u(t) = Q^{-1}B^*\lambda(t)$ for a.e. $t \in [0, T]$, $\bar{u} = Q^{-1}B^*\bar{\lambda}$ and $(\delta y, \delta u, \delta \lambda) := (y - \bar{y}, u - \bar{u}, \lambda - \bar{\lambda})$. Then, if (Λ, B) is exactly controllable and (Λ, C) is exponentially detectable, there exists $\mu > 0$ and a constant $c \geq 0$, both independent of T, such that

$$\begin{split} \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta y\|_{L_2(X)} + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta u\|_{L_2(U)} + \\ + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta \lambda\|_{L_2(X)} &\leq c \left(\|y_0 - \bar{y}\| + \|y_T - \bar{y}\|\right), \\ \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta y\|_{C(X)} + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta u\|_{L_\infty(U)} + \\ + \|\frac{1}{e^{-\mu t} + e^{-\mu(T-t)}} \delta \lambda\|_{C(X)} &\leq c \left(\|y_0 - \bar{y}\| + \|y_T - \bar{y}\|\right). \end{split}$$

Remark 7.3. Similar to Section 5, these results could be extended to the case of unbounded but admissible control, where one again faces the lack of Pontryagin Maximum Principles. Thus, one can only derive results for solutions to the extremal equations instead of the optimal solution, as it was done in Section 5. Additionally, in contrast to the case of no terminal condition, it remains unclear, how the solvability of (7.1), i.e. invertibility of \widetilde{M} can be established, as the approach of Theorem 5.4 fails for initial and terminal conditions on the state.

8. Examples. Finally we provide examples for which the error estimation result for bounded control operators, i.e., Theorem 3.2, the error estimation results for unbounded but admissible control operators, i.e., Theorem 5.8 and the turnpike result Theorem 6.2 hold. First, we consider the interior control of a heat equation and secondly, the boundary control of a wave equation.

Example 8.1. (Interior control of an unstable heat equation) Let $\Omega \subset \mathbb{R}^n$ with smooth boundary and $\omega_c \subset \Omega$ be non-empty. For T > 0, we consider the evolution equation

$$\begin{split} \frac{\partial y}{\partial t} &= (\Delta + c^2 I)y + \chi_{\omega_c} u & \text{in } \Omega \times (0, T) \\ y &= 0 & \text{in } \partial \Omega \times (0, T) \\ y(x, 0) &= y_0 & \text{in } \Omega, \end{split}$$

where χ_{ω_c} is the characteristic function of the control domain ω_c , $U = L_2(\omega_c)$ and $y_0 \in L_2(\Omega)$. Moreover, we consider an observation operator $C = \chi_{\omega_o}$ for non-empty $\omega_o \subset \Omega$. As Δ generates a semigroup on $X = L_2(\Omega)$ with domain $D(\Delta) = H^2(\Omega) \cap H_0^1(\Omega)$, we obtain by classical perturbation results, cf. [9, Theorem 1.3], that $(\Delta + c^2I)$ generates a semigroup with the same domain as Δ . If $c^2 > \lambda_1$, where λ_1 is the smallest eigenvalue of the negative Laplacian, the uncontrolled system is unstable. Defining $B: L_2(\omega_c) \to L_2(\Omega)$ via $Bu := \chi_{\omega_c} u$, the pair (Δ, B) is exactly null controllable in finite time [3, 11], thus stabilizable. Analogously, it follows that (Δ, C) is detectable. As B is bounded on $L_2(X)$, one could apply the sensitivity result Theorem 3.2 or the turnpike result Theorem 6.2 to optimal control problems governed by these operators. For moderate instability parameters c > 0 with respect to the control and observation region, one can also obtain a turnpike and sensitivity result in the W([0,T])-norm with [15, Theorem 3.1, Theorem 5.2] and we refer to [15, Example 3.5] for this matter. However, if c is large, the results of [15] cannot be applied, while the results of the present paper still hold. Concerning the case of terminal constraints it is well known that whenever $\omega_c \neq \Omega$, the system is not exactly controllable and the results of Section 7

can not be applied. This is due to the smoothing effect of the heat equation. For a discussion on this topic, we refer to [37, Chapter 3].

Example 8.2 (Dirichlet control of a wave equation). Second, we provide an example of a hyperbolic PDE with unbounded but admissible control operator presented in [31, Section 10.9]. We consider the model of a vibrating membrane on $\Omega \subset \mathbb{R}^2$, where Ω is a bounded C^2 domain and we can take action through Dirichlet boundary control on a part $\Gamma \subset \partial \Omega$ of the boundary. Moreover, we assume that (Ω, Γ, T) fulfill the so called Geometric Control Condition (GCC), which ensures that all geometric optics have to enter the control domain in a time smaller than T. A consequence of this condition is exact controllability, see [4, 25], and hence stabilizability. We consider the wave equation

$$\begin{split} \frac{\partial^2 w}{\partial t^2} &= \Delta w & \text{on } \Omega \times (0,T) \\ w &= 0 & \text{on } \partial \Omega \setminus \Gamma \times (0,T) \\ w &= u & \text{on } \Gamma \times (0,T) \\ w(x,0) &= f(x), \ \frac{\partial w}{\partial t}(x,0) = g(x) \qquad x \in \Omega, \end{split}$$

where $f \in L_2(\Omega), g \in H^{-1}(\Omega)$ and $u \in L_2(L_2(\Gamma))$. It was shown in [31, Proposition 10.9.1], that one can deduce a corresponding well-posed boundary control system on $X = L_2(\Omega) \times H^{-1}(\Omega)$ with generator $\Lambda = \begin{pmatrix} 0 & I \\ -\Lambda_0 & 0 \end{pmatrix}$, $D(\Lambda) = H_0^1(\Omega) \times L_2(\Omega)$, where Λ_0 is the Dirichlet Laplacian and a control operator B defined by

$$Bv = \begin{pmatrix} 0 \\ \Lambda_0 Dv \end{pmatrix} \quad \forall v \in U = L_2(\Gamma)$$
$$B^* \begin{pmatrix} \varphi \\ \psi \end{pmatrix} = -\frac{\partial}{\partial \nu} (\Lambda_0^{-1} \psi)_{\mid \Gamma} \quad \forall (\varphi, \psi) \in D(\Lambda),$$

where D is the Dirichlet map and $\frac{\partial}{\partial \nu}$ the outward normal derivative. We refer the reader to [31, Section 10.6, Section 10.9] for details. In particular, B is an admissible control operator for the semigroup generated by Λ . Moreover as the GCC is satisfied, the pair (Λ, B) is exponentially stabilizable. If we now consider any observation region $\omega_o \subset \Omega$, such that (Ω, ω_o, T) satisfy the GCC, then the pair (Λ, C) is exponentially detectable. As in this example the control operator is unbounded on X, one cannot apply Theorem 3.2. Yet, Theorem 5.8 and Theorem 6.2 are applicable to the optimal control problems governed by this equation. Finally, as the system is exactly controllable, see [37, Chapter 3], the assumptions of Section 7 are satisfied and the results also hold true if we impose a terminal condition on the state.

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- [1] B. D. Anderson and P. V. Kokotovic, *Optimal control problems over large time intervals*, Automatica, 23 (1987), pp. 355 363.
- [2] J. Ball, Strongly continuous semigroups, weak solutions, and the variation of constants formula, Proceedings of the American Mathematical Society, 63 (1977), pp. 370–373.
- [3] V. Barbu, Controllability of parabolic and Navier-Stokes equations, Scientiae Mathematicae Japonicae, 56 (2002), pp. 143-211.
- [4] C. Bardos, G. Lebeau, and J. Rauch, Sharp sufficient conditions for the observation, control, and stabilization of waves from the boundary, SIAM Journal on Control and Optimization, 30 (1992), pp. 1024–1065.
- [5] T. Breiten and L. Pfeiffer, On the turnpike property and the receding-horizon method for linear-quadratic optimal control problems. 2018, https://arxiv.org/abs/1811.02421.
- [6] J.-M. CORON, Control and nonlinearity, no. 136, American Mathematical Soc., 2007.
- [7] T. DAMM, L. GRÜNE, M. STIELER, AND K. WORTHMANN, An exponential turnpike theorem for dissipative discrete time optimal control problems, SIAM Journal on Control and Optimization, 52 (2014), p. 19351957.
- [8] R. DORFMAN, P. A. SAMUELSON, AND R. M. SOLOW, Linear programming and economic analysis, Courier Corporation, 1958.
- [9] K.-J. ENGEL AND R. NAGEL, One-Parameter Semigroups for Linear Evolution Equations, Springer-Verlag New York, 2000.
- [10] T. FAULWASSER, M. KORDA, C. N. JONES, AND D. BONVIN, On turnpike and dissipativity properties of continuous-time optimal control problems, Automatica, 81 (2017), p. 297304.
- [11] A. FURSIKOV AND O. IMANUVILOV, Controllability of evolution equations, Suhak kangŭirok, Seoul National University, 1996.
- [12] L. GRÜNE, Approximation properties of receding horizon optimal control, Jahresbericht Deutsche Mathematiker-Vereinigung, 118 (2016), pp. 3–37.
- [13] L. GRÜNE AND M. A. MÜLLER, On the relation between strict dissipativity and turnpike properties, Systems & Control Letters, 90 (2016), p. 4553.
- [14] L. GRÜNE AND J. PANNEK, Nonlinear Model Predictive Control: Theory and Algorithms, Springer, 2016.
- [15] L. GRÜNE, M. SCHALLER, AND A. SCHIELA, Sensitivity analysis of optimal control for a class of parabolic PDEs motivated by model predictive control, SIAM Journal on Control and Optimization, 57 (2019), pp. 2753– 2774, https://doi.org/10.1137/18M1223083.
- [16] M. GUGAT AND F. HANTE, On the turnpike phenomenon for optimal boundary control problems with hyperbolic systems, SIAM Journal on Control and Optimization, 57 (2019), pp. 264–289.
- [17] M. GUGAT, E. TRÉLAT, AND E. ZUAZUA, Optimal Neumann control for the 1D wave equation: Finite horizon, infinite horizon, boundary tracking terms and the turnpike property, Systems & Control Letters, 90 (2016), p. 6170
- [18] A. HAURIE, Optimal control on an infinite time horizon: The turnpike approach, Journal of Mathematical Economics, 3 (1976), pp. 81 – 102.
- [19] R. Kress, Linear integral equations, Springer, 2014.
- [20] X. LI AND J. YONG, Optimal Control Theory for Infinite Dimensional Systems, Birkhäuser, 1995.
- [21] J.-L. LIONS, Exact controllability, stabilization and perturbations for distributed systems, SIAM review, 30 (1988), pp. 1–68.
- [22] A. PAZY, Semigroups of Linear Operators and Applications to Partial Differential Equations, Springer, New York, NY, 1983.
- [23] A. PORRETTA AND E. ZUAZUA, Long time versus steady state optimal control, SIAM Journal on Control and Optimization, 51 (2013), pp. 4242–4273.
- [24] A. PORRETTA AND E. ZUAZUA, Remarks on long time versus steady state optimal control, Mathematical Paradigms of Climate Science, (2016), pp. 67–89.
- [25] J. RAUCH, M. TAYLOR, AND R. PHILLIPS, Exponential decay of solutions to hyperbolic equations in bounded domains, Indiana university Mathematics journal, 24 (1974), pp. 79–86.
- [26] D. SALAMON, Infinite-dimensional linear systems with unbounded control and observation: a functional analytic approach, Transactions of the American Mathematical Society, 300 (1987), pp. 383–431.
- [27] D. Salamon, Realization theory in Hilbert space, Mathematical Systems Theory, 21 (1988), pp. 147–164.
- [28] O. Staffans, Well-posed linear systems, vol. 103, Cambridge University Press, 2005.
- [29] E. TRÉLAT, C. ZHANG, AND E. ZUAZUA, Steady-state and periodic exponential turnpike property for optimal control problems in hilbert spaces, SIAM Journal on Control and Optimization, 56 (2018), pp. 1222–1252.
- [30] E. TRÉLAT AND E. ZUAZUA, The turnpike property in finite-dimensional nonlinear optimal control, Journal of Differential Equations, 258 (2015), pp. 81–114.
- [31] M. Tucsnak and G. Weiss, Observation and control for operator semigroups, Springer Science & Business Media, 2009.

- [32] M. TUCSNAK AND G. Weiss, Well-posed systems -the LTI case and beyond, Automatica, 50 (2014), pp. 1757–1779
- [33] G. Weiss, Regular linear systems with feedback, Mathematics of Control, Signals and Systems, 7 (1994), pp. 23–57.
- [34] G. Weiss, O. J. Staffans, and M. Tucsnak, Well-posed linear systems-a survey with emphasis on conservative systems, International Journal of Applied Mathematics and Computer Science, 11 (2001), pp. 7–33.
- [35] M. ZANON, L. GRÜNE, AND M. DIEHL, Periodic optimal control, dissipativity and MPC, IEEE Transactions on Automatic Control, 62 (2017), pp. 2943–2949.
- [36] A. J. Zaslavski, Turnpike theory of continuous-time linear optimal control problems, Springer, 2016.
- [37] E. Zuazua, Controllability and observability of partial differential equations: some results and open problems, Handbook of Differential Equations: Evolutionary Equations, 3 (2007), pp. 527–621.