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# Linear turnpike theorem 

Emmanuel Trélat*


#### Abstract

The turnpike phenomenon stipulates that the solution of an optimal control problem in large time, remains essentially close to a steady-state of the dynamics, itself being the optimal solution of an associated static optimal control problem. Under general assumptions, it is known that not only the optimal state and the optimal control, but also the adjoint state coming from the application of the Pontryagin maximum principle, are exponentially close to a steady-state, except at the beginning and at the end of the time frame. In such results, the turnpike set is a singleton, which is a steady-state.

In this paper, we establish a turnpike result for finite-dimensional optimal control problems in which some of the coordinates evolve in a monotone way, and some others are partial steadystates of the dynamics. We prove that the discrepancy between the optimal trajectory and the turnpike set is then linear, but not exponential: we thus speak of a linear turnpike theorem.


## 1 Introduction and main results

### 1.1 Reminders on the exponential turnpike phenomenon

Let $n, m \in \mathbb{N}^{*}$ and let $T>0$ be arbitrary. We consider a general nonlinear optimal control problem in $\mathbb{R}^{n}$, in fixed final time $T$, without control constraint:

$$
\begin{align*}
& \dot{x}(t)=f(x(t), u(t)), \quad x(0) \in M_{0}, \quad x(T) \in M_{1}  \tag{1}\\
& \min \int_{0}^{T} f^{0}(x(t), u(t)) d t \tag{2}
\end{align*}
$$

where $f: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}, f^{0}: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}, M_{0}$ and $M_{1}$ are subsets of $\mathbb{R}^{n}$, and $u \in$ $L^{\infty}\left([0, T], \mathbb{R}^{m}\right)$ is the control. For simplicity of the exposition, we assume that there exists $T_{0}>0$ such that, for every $T \geqslant T_{0}$, there exists a unique optimal solution of (1)-(2), denoted by $(x(\cdot), u(\cdot))$ (sufficient conditions ensuring existence are standard, see, e.g., [25]; uniqueness is ensured under differentiability properties of the value function and thus is generic in some sense). Note that $(x(\cdot), u(\cdot))$ depends on $T$ but we do not denote it by $\left(x_{T}(\cdot), u_{T}(\cdot)\right)$ to keep readability. By the Pontryagin maximum principle (see [19]), there exist $p^{0} \leqslant 0$ and an absolutely continuous mapping $p_{x}(\cdot):[0, T] \rightarrow \mathbb{R}^{n}$ (called adjoint vector), satisfying $\left(p_{x}(\cdot), p^{0}\right) \neq(0,0)$, such that

$$
\begin{align*}
& \dot{x}(t)=\frac{\partial H}{\partial p_{x}}\left(x(t), p_{x}(t), p^{0}, u(t)\right), \quad \dot{p}_{x}(t)=-\frac{\partial H}{\partial x}\left(x(t), p_{x}(t), p^{0}, u(t)\right) \\
& \frac{\partial H}{\partial u}\left(x(t), p_{x}(t), p^{0}, u(t)\right)=0 \tag{3}
\end{align*}
$$

[^0]for almost every $t \in[0, T]$, where $H\left(x, p_{x}, p^{0}, u\right)=\left\langle p_{x}, f(x, u)\right\rangle+p^{0} f^{0}(x, u)$ is the Hamiltonian of the optimal control problem (1)-(2) (here and throughout, $\langle\cdot, \cdot\rangle$ is the Euclidean scalar product in $\mathbb{R}^{n}$, and $\|\cdot\|$ is the corresponding norm). We have moreover the transversality conditions on the adjoint vector at initial and final time:
\[

$$
\begin{equation*}
p_{x}(0) \perp T_{x(0)} M_{0}, \quad p_{x}(T) \perp T_{x(T)} M_{1} \tag{4}
\end{equation*}
$$

\]

provided the tangent space $T_{x(0)} M_{0}$ to $M_{0}$ at $x(0)$ and the tangent space $T_{x(T)} M_{1}$ to $M_{1}$ at $x(T)$ exist (these conditions are empty when the initial and final points are fixed in the optimal control problem). We assume that the extremal lift $\left(x(\cdot), p_{x}(\cdot), p^{0}, u(\cdot)\right)$ is not abnormal and thus that one can fix $p^{0}=-1$. We also assume, for the simplicity of the exposition, that the extremal lift is unique (these two conditions are generic, see [6, 7]).

The turnpike phenomenon stipulates that, when the final time $T$ is large, the optimal solution $\left(x(\cdot), p_{x}(\cdot), u(\cdot)\right)$ remains, along the interval of time $[0, T]$ and except around the initial time $t=0$ and around the final time $t=T$, "essentially close" to some static point $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n} \times \mathbb{R}^{m}$, i.e., roughly speaking,

$$
\begin{equation*}
x(t) \simeq \bar{x}, \quad p_{x}(t) \simeq \bar{p}_{x}, \quad u(t) \simeq \bar{u} \quad \forall t \in[\eta, T-\eta] \tag{5}
\end{equation*}
$$

(for some $\eta>0$ ) and moreover, $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right)$ is itself the solution of a minimization problem, called static optimal control problem, which is:

$$
\begin{equation*}
\min \left\{f^{0}(x, u) \mid(x, u) \in \mathbb{R}^{n} \times \mathbb{R}^{m}, f(x, u)=0\right\} \tag{6}
\end{equation*}
$$

i.e., in other words, the problem of minimizing the instantaneous cost $f^{0}(x, u)$ over all possibles equilibrium points of the dynamics $(f(x, u)=0)$. Note that, by the Lagrange multiplier rule, expressed in a Hamiltonian form, applied to the static optimal control problem, we have

$$
\begin{equation*}
\frac{\partial H}{\partial p_{x}}\left(\bar{x}, \bar{p}_{x}, p^{0}, \bar{u}\right)=0, \quad \frac{\partial H}{\partial x}\left(\bar{x}, \bar{p}_{x}, p^{0}, \bar{u}\right)=0, \quad \frac{\partial H}{\partial u}\left(\bar{x}, \bar{p}_{x}, p^{0}, \bar{u}\right)=0 \tag{7}
\end{equation*}
$$

Here we assume as well that the minimizer $(\bar{x}, \bar{u})$ is unique, and also that the quadruple $\left(\bar{x}, \bar{p}_{x}, p^{0}, \bar{u}\right)$ is not abnormal, so that one can take $p^{0}=-1$ (usual qualification condition in optimization theory) and that the Lagrange multiplier is unique (the latter two conditions are generic).

We note then that $\left(\bar{x}, \bar{p}_{x},-1, \bar{u}\right)$ is an equilibrium point of (3). Hence, the informal property (5) stipulates that the optimal triple $\left(x(\cdot), p_{x}(\cdot), u(\cdot)\right)$, solution of (1)-(2)-(3) remains essentially close to an equilibrium point of the dynamics (3), this equilibrium being itself an optimal solution of the static problem (6). This is the turnpike phenomenon.

The intuition explaining why such a property is to be expected is the following. Since $T$ is assumed to be large, we set $\varepsilon=\frac{1}{T}$, which is considered as a small parameter. Now, we make a time reparametrization, by setting $s=\varepsilon t$, so that, when $t$ ranges over $[0, T], s$ ranges over $[0,1]$. Setting $\mathrm{x}_{\varepsilon}(s)=x(t), \mathrm{p}_{\varepsilon}(s)=p_{x}(t)$ and $\mathrm{u}_{\varepsilon}(s)=u(t)$, we have, still in an informal way,

$$
\begin{align*}
& \frac{\partial H}{\partial p_{x}}\left(\mathrm{x}_{\varepsilon}(s), \mathrm{p}_{\varepsilon}(s),-1, \mathrm{u}_{\varepsilon}(s)\right)=\varepsilon \mathrm{x}_{\varepsilon}^{\prime}(s) \simeq 0 \\
& \frac{\partial H}{\partial x}\left(\mathrm{x}_{\varepsilon}(s), \mathrm{p}_{\varepsilon}(s),-1, \mathrm{u}_{\varepsilon}(s)\right)=-\varepsilon \mathrm{p}_{\varepsilon}^{\prime}(s) \simeq 0  \tag{8}\\
& \frac{\partial H}{\partial u}\left(\mathrm{x}_{\varepsilon}(s), \mathrm{p}_{\varepsilon}(s),-1, \mathrm{u}_{\varepsilon}(s)\right)=0
\end{align*}
$$

which gives credibility to (5), at least when $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right)$ is the unique solution of (7) (and under nondegenerate assumptions on the system of equations (7)). But this is not a proof and this
remains completely informal: indeed, to justify (8) one would need that the derivatives $\mathrm{x}_{\varepsilon}^{\prime}(s)$ and $\mathrm{p}_{\varepsilon}^{\prime}(s)$ be bounded (except around $t=0$ and $t=T$ ), uniformly with respect to $\varepsilon$, i.e., that the derivatives $T \dot{x}(t)$ and $T \dot{p}_{x}(t)$ be bounded (except around $t=0$ and $t=T$ ), uniformly with respect to $T$, which is precisely the hard part in establishing a turnpike result.

The main result of [29] states an exponential turnpike property, which is local (see Remark 1). Hereafter, we combine this result with a result of [27] in order to obtain a global exponential turnpike property. We use the notion of strict dissipativity, introduced in [31] and already used in [11, $14,15,16]$ to derive turnpike properties. We recall that the family of optimal control problems (1)-(2), indexed by $T>0$, is strictly dissipative at the optimal static point ( $\bar{y}, \bar{u}$ ) with respect to the supply rate function

$$
\begin{equation*}
w(x, u)=f^{0}(x, u)-f^{0}(\bar{x}, \bar{u}) \tag{9}
\end{equation*}
$$

if there exists a locally bounded function $S: \mathbb{R}^{n} \rightarrow \mathbb{R}$, called storage function, such that

$$
\begin{equation*}
S(x(0))+\int_{0}^{T} w(x(t), u(t)) d t \geqslant S(x(T))+\frac{1}{T} \int_{0}^{T} \alpha(\|x(t)-\bar{x}\|,\|u(t)-\bar{u}\|) d t \tag{10}
\end{equation*}
$$

(strict dissipation inequality) for every $t \in[0, T]$, for every $T>0$, for any pair $(x(\cdot), u(\cdot))$ solution of (1), for some function $\alpha$ of class $\mathcal{K}$, i.e., $\alpha:[0,+\infty) \rightarrow[0,+\infty)$ continuous, increasing, and $\alpha(0)=0$. Comments on dissipativity can be found in [4, 27, 31].

We use hereafter the shorter notations

$$
H_{\star \#}=\frac{\partial^{2} H}{\partial \star \partial \#}(\bar{x}, \bar{u})
$$

where $\star$ and \# will be replaced with such or such variable. We set

$$
\bar{A}_{1}=H_{p_{x} x}=\frac{\partial f}{\partial x}(\bar{x}, \bar{u}), \quad \bar{B}_{1}=H_{p_{x} u}=\frac{\partial f}{\partial u}(\bar{x}, \bar{u}), \quad U=-H_{u u}
$$

and, since we assume hereafter that $U$ is invertible,

$$
A_{1}=\bar{A}_{1}+\bar{B}_{1} U^{-1} H_{u x}, \quad W=-H_{x x}-H_{x u} U^{-1} H_{u x}
$$

Theorem 1. We make the following assumptions:
Assumptions of global nature:
(i) There exist $R>0$ and $T_{0}>0$ such that, for every $T \geqslant T_{0}$, there exists a unique optimal triple $\left(x(\cdot), p_{x}(\cdot), u(\cdot)\right)$ solution of (1)-(2)-(3)-(4) (assumed to be normal). Moreover, $\|x(t)\|+$ $\left\|p_{x}(t)\right\|+\|u(t)\| \leqslant R$ for every $t \in[0, T]$.
(ii) There exists a unique optimal triple $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right)$ solution of (6)-(7) (assumed to be normal).
(iii) The family of optimal control problems (1)-(2), indexed by $T>0$, is strictly dissipative at the optimal static point $(\bar{y}, \bar{u})$ with respect to the supply rate function $w$ defined by (9), with a storage function $S$ (satisfying (10)) that is bounded on $M_{0}$ and $M_{1}$.
(iv) There exist $\delta>0$ and a trajectory $x_{1}(\cdot)$ (resp., $x_{2}(\cdot)$ ), solution of (1) on some time interval $\left[0, t_{1}\right]$, with $t_{1} \leqslant \delta$ (resp., $\left[0, t_{2}\right]$ with $t_{2} \leqslant \delta$ ), such that $x(0) \in M_{0}$ and $x\left(t_{1}\right)=\bar{x}$ (resp., such that $x(0)=\bar{x}$ and $\left.x\left(t_{2}\right) \in M_{1}\right)$.

Assumptions of local nature:
(iv) $U=-H_{u u}=-\frac{\partial^{2} H}{\partial u^{2}}(\bar{x}, \bar{u})$ is a positive definite symmetric matrix.
(v) $W$ is a positive definite symmetric matrix.
(vi) The pair $\left(\bar{A}_{1}, \bar{B}_{1}\right)$ satisfies the Kalman condition ${ }^{1}$, i.e., $\operatorname{rank}\left(\bar{B}_{1}, \bar{A}_{1} \bar{B}_{1}, \ldots, \bar{A}_{1}^{n-1} \bar{B}_{1}\right)=n$. In other words, the linearized system at $(\bar{x}, \bar{u})$ is controllable.

Then, there exist $C>0$ and $\nu>0$ such that if $T \geqslant T_{0}$ then

$$
\begin{equation*}
\|x(t)-\bar{x}\|+\left\|p_{x}(t)-\bar{p}_{x}\right\|+\|u(t)-\bar{u}\| \leqslant C\left(e^{-\nu t}+e^{-\nu(T-t)}\right) \quad \forall t \in[0, T] \tag{11}
\end{equation*}
$$

In other words, except around $t=0$ and $t=T$, the optimal triple $\left(x(\cdot), p_{x}(\cdot), u(\cdot)\right)$ solution of (1)-(2)-(3)-(4) is exponentially close to the optimal triple $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right)$ solution of (6)-(7). The constants $C$ and $\nu$ can even be made explicit by solving a Riccati equation (see [29] for details).

The local assumptions (iv)-(v)-(vi) done in this theorem have been chiefly discussed and commented in [29]. The assumption (iv) is a global controllability assumption: there exists a trajectory steering in finite time the control system (1) from $M_{0}$ to the turnpike $\bar{x}$, and from the turnpike $\bar{x}$ to the final set $M_{1}$. The assumption (i) on the uniqueness of the solution of the optimal control problem and normality of the extremal triple is actually generic, in some sense: indeed it is well known that it is satisfied as soon as the value function corresponding to the optimal control problem is continuous at the terminal points under consideration. But, such comments are completely general in optimal control theory and we thus do not comment more on the generality of our assumptions.

Remark 1. The above statement does not appear as such in the literature. Compared with [29] where the main result is established under the assumptions of local nature only (i.e., (iv)-(v)(vi)), we have added here assumptions of global nature (i.e., (i)-(ii)-(iii)-(iv)): compactness and uniqueness of minimizers, strict dissipativity, global controllability to and from the turnpike set, and then we obtain a global turnpike property.

When the minimizers are not unique or when the system is not dissipative, we obtain only a local turnpike property, in the sense that the turnpike estimate (11) can only be expected to be satisfied in a neighborhood of an extremal steady-state $\left(\bar{x}, \bar{p}_{x}, \bar{u}\right)$ (see [29]). This is so, because the dynamics is nonlinear. In dynamical systems theory, when an equilibrium is non degenerate, the solutions of the nonlinear system resemble the solutions of the linearized system at the equilibrium, only locally around the equilibrium. The above turnpike theorem is of this nature.

It is worth noting that the exponential turnpike estimate (11) is global for non-degenerate linear-quadratic optimal control problems (by strict convexity), without having to assume (i)-(ii)(iii).

When the static problem has several local minimizers, then one has local turnpike properties. When the static problem has, for instance, two global minimizers, then there is a region of the state space in which globally optimal solutions of the optimal control problems have a turnpike property corresponding to the first global minimizer, and another region in which the turnpike property corresponds to the second global minimizer. We refer to Appendix A for a discussion on these topics and for some numerical simulations.

The turnpike result of [29] has been extended to the infinite-dimensional setting in [27, 28, 10, 17]. One can find many turnpike results in the literature (see [5, 11, 16, 15, 20, 21, 32, 33], and see historical references in [29]), but the specificity of the above result is that the turnpike property is established as well for the adjoint vector: as explained in [29], this is particularly important in view of deriving appropriate initializations in numerical computation methods (direct methods or shooting methods, see [26]).

[^1]Remark 2. It is proved in [27] that "strong duality implies dissipativity". To be more precise, we recall that the static problem (6) has the strong duality property if $(\bar{x}, \bar{u})$ minimizes the Lagrangian function $L\left(\cdot, \cdot, \bar{p}_{x}\right): \mathbb{R}^{n} \times \mathbb{R}^{m} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ defined by $L\left(x, u, \bar{p}_{x}\right)=f^{0}(x, u)-\left\langle\bar{p}_{x}, f(x, u)\right\rangle$ (the sign minus is due to the fact that we took $p^{0}=-1$ for the Lagrange multiplier associated with the cost). This notion is well known in classical optimization theory, in relation with primal and dual problems. It is proved in [27, Theorem 3] that the strong duality property implies the dissipativity property, with the storage function $S(x)=\left\langle\bar{p}_{x}, x\right\rangle$.

Note that strong duality is in some sense an infinitesimal version of the dissipation inequality, at least when the storage function is continuously differentiable (see also [11, 13]).

Until now, we have discussed a turnpike phenomenon around an equilibrium point of the dynamics. In the next section, we state a turnpike theorem for systems having some coordinates that evolve in a monotone way and thus have no equilibrium point.

### 1.2 A new result: linear turnpike phenomenon

### 1.2.1 Framework

We keep the notations introduced in Section 1.1, and we add to the control system (1)-(2) p additional differential equations, for some $p \in \mathbb{N}^{*}$, as follows. We consider the general optimal control problem in $\mathbb{R}^{n+p}$ in fixed final time $T>0$, without control constraint:

$$
\begin{array}{lrl}
\dot{x}(t)=f(x(t), u(t)), & x(0) \in M_{0}, & x(T) \in M_{1} \\
\dot{y}(t)=g(x(t), u(t)), & y(0)=y_{0}, & y(T)=y_{1}^{T} \\
\min \int_{0}^{T} f^{0}(x(t), u(t)) d t & & \tag{14}
\end{array}
$$

where $y_{0}, y_{1}^{T} \in \mathbb{R}^{p}$ are fixed and $g: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{p}$ is of class $C^{2}$ (the rest of the assumptions is unchanged, with respect to Section 1.1).

Compared with (1)-(2), we have added the $p$ differential equations $\dot{y}(t)=g(x(t), u(t))$, where $g$ does not depend on the new coordinate $y \in \mathbb{R}^{p}$. The instantaneous cost function $f^{0}$ also, still does not depend on $y$. This is important to derive the result hereafter.

As in the previous section, the final time $T$ is fixed in the optimal control problem (12)-(13)-(14) and the linear turnpike property will describe the behavior of optimal solutions when $T$ is large.

In the optimal control problem (12)-(13)-(14), the data $M_{0}, M_{1}$ and $y_{0}$ do not depend on $T$, but the final condition $y_{1}^{T}$ depends on $T$ : we have thus put a superscript $T$ in order to underline this important dependence. We have indeed in mind to consider cases where $\left\|y_{1}^{T}\right\|$ is large as $T$ is large: in such cases, we expect that the components of $y(t)$ should evolve in a monotone way, i.e., $g(x(t), u(t)) \neq 0$ along the optimal trajectory: in such a case, we have no equilibrium in the $y$ components.

We are going to derive a turnpike result in which the "turnpike set" consists of a partial equilibrium, i.e., the turnpike is $(\bar{x}, \bar{y}(t), \bar{u}) \in \mathbb{R}^{n} \times \mathbb{R}^{p} \times \mathbb{R}^{m}$ where $f(\bar{x}, \bar{u})=0$, while its $y$ component $t \mapsto \bar{y}(t)$ evolves in a monotone way.

Remark 3. Following [9], it is worth noting that such problems (12)-(13)-(14) frequently appear in applications, for instance in mechanical problems where, denoting by $\mathrm{x}(t)$ a position variable, by $\mathrm{v}(t)=\dot{\mathrm{x}}(t)$ the speed variable, one has a system of the form

$$
\dot{\mathrm{x}}(t)=\mathrm{v}(t), \quad \dot{\mathrm{v}}(t)=f(\mathrm{v}(t), \mathrm{u}(t))
$$

The result that we are going to derive allows one to obtain a turnpike phenomenon along a trajectory $\mathrm{x}(t)=\overline{\mathrm{v}} t$, where $\overline{\mathrm{v}} \in \mathbb{R}^{n}$ is such that $f(\overline{\mathrm{v}}, \overline{\mathrm{u}})=0$. Here, $(\overline{\mathrm{v}}, \overline{\mathrm{u}})$ is an equilibrium of $f$, but
$(\overline{\mathrm{v}} t, \overline{\mathrm{v}}, \overline{\mathrm{u}})$ is not an equilibrium of the complete dynamics, and the component $x(t)=\bar{v} t$ evolves in a monotone way. In [9], this phenomenon is called a "velocity turnpike". Because of the estimate that we are going to establish in Theorem 2 hereafter, we rather speak of a "linear turnpike" phenomenon.

Note that, in (12)-(13)-(14), we have fixed the initial and final conditions on $y(t)$. Indeed, otherwise, if either $y(0)$ or $y(T)$ is let free, then there is nothing new, because $f$ and $f^{0}$ do not depend on $y$ and then (12)-(13)-(14) is then completely equivalent to (1)-(2), with the coordinate $y(t)$ evolving without having any impact on the optimal control problem (1)-(2). But here, we impose $y(0)=y_{0}$ and $y(T)=y_{1}^{T}$, which creates additional constraints with respect to (1)-(2), and then the turnpike result stated in Theorem 1 cannot be applied.

### 1.2.2 A motivating example

In order to shed light on what we can expect, let us consider the very simple example in $\mathbb{R}^{2}$ :

$$
\begin{array}{lll}
\dot{x}(t)=u(t), & x(0)=1, & x(T)=2 \\
\dot{y}(t)=x(t), & y(0)=0, & y(T)=y_{1}^{T}  \tag{15}\\
\min \int_{0}^{T}\left(x(t)^{2}+u(t)^{2}\right) d t & &
\end{array}
$$

on which the optimal trajectory can be computed explicitly. Here, $y_{1}^{T} \in \mathbb{R}$ is a smooth function of $T$ that is sublinear at infinity. For instance, one can take $y_{1}^{T}=\alpha T$ or $y_{1}^{T}=\alpha T \sin (T)$ for some $\alpha \in \mathbb{R}$.

In a first step, if we ignore the dynamics in $y$ in (15), then we are in the framework of Theorem 1. The static optimal control problem consists of minimizing $x^{2}+u^{2}$ under the constraint $u=0$, which gives the optimal solution $\bar{x}=0, \bar{u}=0$, and the Lagrange multiplier $\bar{p}_{x}=0$ (which is normal; here, we take $p^{0}=-1 / 2$ ). By Theorem 1 , we obtain $|x(t)|+\left|p_{x}(t)\right|+|u(t)| \leqslant C\left(e^{-\nu t}+e^{-\nu(T-t)}\right)$ for every $t \in[0, T]$. In other words, roughly, $x(t) \simeq 0$ and $u(t) \simeq 0$, except around $t=0$ and $t=T$.

Now, let us take into account the dynamics in $y$ and the terminal conditions on $y$. We observe that we are not anymore in the framework of Theorem 1, because there is no term in $y(t)^{2}$ in the cost and therefore the matrix $W$ is not positive definite. We define the "static" optimal control problem as

$$
\begin{equation*}
\min \left\{x^{2}+u^{2} \quad \mid \quad(x, u) \in \mathbb{R} \times \mathbb{R}, \quad u=0, \quad T x=y_{1}^{T}\right\} \tag{16}
\end{equation*}
$$

i.e., with respect to the previous static problem, we add the terminal constraints coming from the dynamics in $y$ : the unique solution of the Cauchy problem $\dot{y}(t)=x$ (with $x$ constant), $y(0)=0$, must satisfy $y(T)=y_{1}^{T}$.

It is easy to check that the (unique) optimal solution of (16) is $\bar{x}=\frac{y_{1}^{T}}{T}, \bar{u}=0$, with Lagrange multipliers $\bar{p}_{x}=0$ associated with the constraint $u=0$, and $\bar{p}_{y}=\frac{y_{1}^{T}}{T}$ associated with the constraint $x=\frac{y_{1}^{T}}{T}$ (which are normal; here, we take $p^{0}=-1 / 2$ ). The corresponding trajectory in $y$ is $\bar{y}(t)=\frac{y_{1}^{T}}{T} t$ (and we have $\bar{y}(T)=y_{1}^{T}$ ). For this example, we expect that, if $T>0$ is large, then $x(t) \simeq 0, u(t) \simeq 0$ and $y(t) \simeq \bar{y}(t)$ for $t \in[0, T]$, where the order of this approximation is to be determined, what we do hereafter by explicit calculations.

Let us apply the Pontryagin maximum principle to (15). The Hamiltonian is $H=p_{x} u+p_{y} x-$ $\frac{1}{2}\left(x^{2}+u^{2}\right)$ (there is no abnormal minimizer). Then $u(t)=p_{x}(t)$ and the extremal system is

$$
\dot{x}(t)=p_{x}(t), \quad \dot{y}(t)=x(t), \quad \dot{p}_{x}(t)=x(t)-p_{y}
$$

with $p_{y}$ constant. We have then

$$
\begin{aligned}
\frac{d}{d t}\binom{x(t)-p_{y}}{p_{x}(t)} & =\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)\binom{x(t)-p_{y}}{p_{x}(t)} \\
\frac{d}{d t} y(t) & =x(t)
\end{aligned}
$$

In particular, the dynamics in $\left(x-p_{y}, p_{x}\right)$ is hyperbolic: setting $v(t)=x(t)-p_{y}+p_{x}(t)$ and $w(t)=x(t)-p_{y}-p_{x}(t)$, we have

$$
\dot{v}(t)=v(t), \quad \dot{w}(t)=-w(t)
$$

with terminal conditions $v(0)+w(0)=2-2 p_{y}, v(T)+w(T)=4-2 p_{y}$. The solutions of this differential system are given by the familiar phase portrait around a saddle point, and we infer that

$$
|v(t)|+|w(t)|=\mathrm{O}\left(e^{-t}+e^{-(T-t)}\right)+\mathrm{O}\left(\left|p_{y}\right|\left(e^{-t}+e^{-(T-t)}\right)\right)
$$

for every $t \in[0, T]$, where the $\mathrm{O}(\cdot)$ is uniform with respect to $T$, and thus

$$
\left|x(t)-p_{y}\right|+\left|p_{x}(t)\right|=\mathrm{O}\left(e^{-t}+e^{-(T-t)}\right)+\mathrm{O}\left(\left|p_{y}\right|\left(e^{-t}+e^{-(T-t)}\right)\right)
$$

In particular,

$$
x(t)=p_{y}+\mathrm{O}\left(e^{-t}+e^{-(T-t)}\right)+\mathrm{O}\left(\left|p_{y}\right|\left(e^{-t}+e^{-(T-t)}\right)\right)
$$

and, integrating $\dot{y}(t)=x(t)$ we obtain $y(t)=p_{y} t+\mathrm{O}(1)+\mathrm{O}\left(\left|p_{y}\right|\right)$ for every $t \in[0, T]$. Since $y(T)=y_{1}^{T}$, we obtain that

$$
p_{y}=\frac{y_{1}^{T}}{T}+\mathrm{O}\left(\frac{1}{T}\right)+\mathrm{O}\left(\frac{\left|p_{y}\right|}{T}\right)
$$

and since $y_{1}^{T}$ is sublinear at infinity, we conclude that

$$
p_{y}=\bar{p}_{y}+\mathrm{O}\left(\frac{1}{T}\right)
$$

Therefore

$$
x(t)=\frac{y_{1}^{T}}{T}+\mathrm{O}\left(\frac{1}{T}+e^{-t}+e^{-(T-t)}\right)=\bar{x}+\mathrm{O}\left(\frac{1}{T}+e^{-t}+e^{-(T-t)}\right)
$$

and

$$
y(t)=\frac{y_{1}^{T}}{T} t+\mathrm{O}(1)=\bar{y}(t)+\mathrm{O}(1)
$$

for every $t \in[0, T]$. In turn, we have also $p_{x}(t)=\mathrm{O}\left(e^{-t}+e^{-(T-t)}\right)$.
We conclude that there exists $C>0$ such that

$$
\begin{align*}
& |x(t)-\bar{x}| \leqslant C\left(\frac{1}{T}+e^{-t}+e^{-(T-t)}\right)  \tag{17}\\
& |u(t)-\bar{u}|+\left|p_{x}(t)-\bar{p}_{x}\right| \leqslant C\left(e^{-t}+e^{-(T-t)}\right) \\
& |y(t)-\bar{y}(t)| \leqslant C \\
& \left|p_{y}-\bar{p}_{y}\right| \leqslant \frac{C}{T}
\end{align*}
$$

for every $t \in[0, T]$. With respect to the conclusion of Theorem 1 , we see on this example that we do not have an exponential turnpike phenomenon for the component $x(t)$ : the above inequality is weaker; it says that, except around $t=0$ and $t=T, x(t)$ is bounded by $\frac{1}{T}$ (which is small as $T$ is large), while the component $y(t)$ remains in a bounded neighborhood of $\bar{y}(t)$.

In what follows, we are going to generalize the above example.

### 1.2.3 Main result

As in Section 1.1, we assume that there exists a unique optimal solution $(x(\cdot), y(\cdot), u(\cdot))$ of (12)-(13)-(14). With respect to (1)-(2), the Hamiltonian of (12)-(13)-(14) is now

$$
H\left(x, y, p_{x}, p_{y}, p^{0}, u\right)=\left\langle p_{x}, f(x, u)\right\rangle+\left\langle p_{y}, g(x, u)\right\rangle+p^{0} f^{0}(x, u)
$$

Assuming as well that the optimal solution has a unique extremal lift, which is not abnormal (so that we take $p^{0}=-1$ ), the application of the Pontryagin maximum principle leads to the existence of absolutely continuous mappings $p_{x}(\cdot):[0, T] \rightarrow \mathbb{R}^{n}$ and $p_{y}(\cdot):[0, T] \rightarrow \mathbb{R}^{p}$ such that

$$
\begin{align*}
\dot{x}(t) & =\frac{\partial H}{\partial p_{x}}\left(x(t), y(t), p_{x}(t), p_{y}(t),-1, u(t)\right)=f(x(t), u(t)) \\
\dot{y}(t) & =\frac{\partial H}{\partial p_{y}}\left(x(t), y(t), p_{x}(t), p_{y}(t),-1, u(t)\right)=g(x(t), u(t)) \\
\dot{p}_{x}(t) & =-\frac{\partial H}{\partial x}\left(x(t), y(t), p_{x}(t), p_{y}(t),-1, u(t)\right)  \tag{18}\\
& =-\frac{\partial f}{\partial x}(x(t), u(t))^{\top} p_{x}(t)-\frac{\partial g}{\partial x}(x(t), u(t))^{\top} p_{y}(t)+\frac{\partial f^{0}}{\partial x}(x(t), u(t))^{\top} \\
\dot{p}_{y}(t) & =-\frac{\partial H}{\partial y}\left(x(t), y(t), p_{x}(t), p_{y}(t),-1, u(t)\right)=0
\end{align*}
$$

and

$$
\begin{align*}
\frac{\partial H}{\partial u}\left(x(t), y(t), p_{x}(t)\right. & \left., p_{y}(t),-1, u(t)\right) \\
& =\frac{\partial f}{\partial u}(x(t), u(t))^{\top} p_{x}(t)+\frac{\partial g}{\partial u}(x(t), u(t))^{\top} p_{y}(t)-\frac{\partial f^{0}}{\partial x}(x(t), u(t))^{\top}=0 \tag{19}
\end{align*}
$$

for almost every $t \in[0, T]$. Note that $p_{y}(t)=p_{y} \in \mathbb{R}^{p}$ is constant.
With respect to (6), the "static" optimal control problem is now the following one:

$$
\begin{align*}
\min \left\{f^{0}(x, u) \mid(x, u) \in \mathbb{R}^{n} \times \mathbb{R}^{m}, \quad\right. & f(x, u)=0 \\
& \left.\dot{y}(t)=g(x, u), \quad y(0)=y_{0}, \quad y(T)=y_{1}^{T}\right\} \tag{20}
\end{align*}
$$

where we have added the constraint coming from the dynamics and terminal conditions in $y$. As before, we assume that there exists a unique optimal solution $(\bar{x}, \bar{y}, \bar{u})$ of (20) and that it has a unique extremal (Lagrange multiplier) lift, which is not abnormal. Integrating the differential equation in (20), noting that $(x, u)$ is constant, the "static" optimal control problem (20) is equivalent to

$$
\begin{equation*}
\min \left\{f^{0}(x, u) \mid(x, u) \in \mathbb{R}^{n} \times \mathbb{R}^{m}, \quad f(x, u)=0, \quad g(x, u)=\frac{y_{1}^{T}-y_{0}}{T}\right\} \tag{21}
\end{equation*}
$$

which is a classical optimization problem under equality constraints. The Lagrange multiplier rule, applied to the solution $(\bar{x}, \bar{y}, \bar{u})$ of this optimization problem, implies (since it is assumed not to be abnormal) that $d f^{0}(\bar{x}, \bar{u})$ can be expressed linearly in function of $d f(\bar{x}, \bar{u})$ and $d g(\bar{x}, \bar{u})$. In the next lemma, we write these conditions in an equivalent Hamiltonian form, which is more complicated but is exactly devised to be comparable with (18).

Lemma 1. There exist $\bar{p}_{x} \in \mathbb{R}^{n}$ and $\bar{p}_{y} \in \mathbb{R}^{p}$ such that, for every $t \in[0, T]$,

$$
\begin{aligned}
0 & =\frac{\partial H}{\partial p_{x}}\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y},-1, \bar{u}\right)=f(\bar{x}, \bar{u}) \\
\dot{\bar{y}}(t) & =\frac{\partial H}{\partial p_{y}}\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y},-1, \bar{u}\right)=g(\bar{x}, \bar{u}), \quad \bar{y}(0)=y_{0}, \quad \bar{y}(T)=y_{1}^{T} \\
0 & =-\frac{\partial H}{\partial x}\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y},-1, \bar{u}\right)=-\frac{\partial f}{\partial x}(\bar{x}, \bar{u})^{\top} \bar{p}_{x}-\frac{\partial g}{\partial x}(\bar{x}, \bar{u})^{\top} \bar{p}_{y}+\frac{\partial f^{0}}{\partial x}(\bar{x}, \bar{u}) \\
0 & =-\frac{\partial H}{\partial y}\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y},-1, \bar{u}\right)
\end{aligned}
$$

and

$$
\frac{\partial H}{\partial u}\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y},-1, \bar{u}\right)=\frac{\partial f}{\partial u}(\bar{x}, \bar{u})^{\top} \bar{p}_{x}+\frac{\partial g}{\partial u}(\bar{x}, \bar{u})^{\top} \bar{p}_{y}-\frac{\partial f^{0}}{\partial u}(\bar{x}, \bar{u})=0 .
$$

This lemma follows by applying the Lagrange multiplier rule to the optimization problem (21) (see also Lemma 2 in Section 2).

We observe that $\bar{y}(t)=y_{0}+\operatorname{tg}(\bar{x}, \bar{u})$ and therefore $y_{0}+T g(\bar{x}, \bar{u})=y_{1}^{T}$. When $g(\bar{x}, \bar{u}) \neq 0$, this can only be if the norm of $y_{1}^{T}-y_{0}$ is large, as $T$ is large. In the proof of the turnpike theorem hereafter, we are going to consider $T$ as a parameter, tending to $+\infty$. This is why we have put a superscript $T$ in $y_{1}^{T}$ : we indeed want to underline that, when $T$ becomes larger, it is required that $y_{1}^{T}-y_{0}$ becomes larger as well (otherwise, the controllability problem may have no solution).

We use the same notations as before, and additionally, we set

$$
\begin{gathered}
\bar{A}_{2}=H_{p_{y} x}=\frac{\partial g}{\partial x}(\bar{x}, \bar{u}), \quad \bar{B}_{2}=H_{p_{y} u}=\frac{\partial g}{\partial u}(\bar{x}, \bar{u}), \quad A_{2}=\bar{A}_{2}+\bar{B}_{2} U^{-1} H_{u x}, \\
\bar{A}=\binom{\bar{A}_{1}}{\bar{A}_{2}}=\binom{\frac{\partial f}{\partial x}(\bar{x}, \bar{u})}{\frac{\partial g}{\partial x}(\bar{x}, \bar{u})}, \quad \bar{B}=\binom{\bar{B}_{1}}{\bar{B}_{2}}=\binom{\frac{\partial f}{\partial u}(\bar{x}, \bar{u})}{\frac{\partial g}{\partial u}(\bar{x}, \bar{u})} .
\end{gathered}
$$

In our main result hereafter, we make all assumptions done in Theorem 1 and we adapt them in an obvious way to the present context where the turnpike set is not anymore the singleton $\{\bar{x}\}$, but is now the trajectory $t \mapsto(\bar{x}, \bar{y}(t))$. For instance, the controllability assumption (iv) now stipulates that there exists a trajectory steering in finite time some point of the initial set $M_{0}$ to some point of the turnpike trajectory, and that any point of the turnpike trajectory can be steered in finite time to the final set $M_{1}$.

Theorem 2. In addition to the assumptions done in Theorem 1, we assume that:

- $\operatorname{ker}\left(\bar{A}^{\top}\right) \cap \operatorname{ker}\left(\bar{B}^{\top}\right)=\{0\} ;^{2}$
- there exist $C_{1}>0$ and $T_{0}>0$ such that $\left\|y_{1}^{T}\right\| \leqslant C_{1} T$ for every $T \geqslant T_{0}$.

Then, there exist $C>0$ and $\nu>0$ such that if $T \geqslant T_{0}$ then

$$
\begin{align*}
& \|x(t)-\bar{x}\|+\left\|p_{x}(t)-\bar{p}_{x}\right\|+\|u(t)-\bar{u}\| \leqslant C\left(\frac{1}{T}+e^{-\nu t}+e^{-\nu(T-t)}\right)  \tag{22}\\
& \|y(t)-\bar{y}(t)\| \leqslant C, \quad\left\|p_{y}-\bar{p}_{y}\right\| \leqslant \frac{C}{T}
\end{align*}
$$

for every $t \in[0, T]$.

[^2]The first estimate in (22) states that, except around $t=0$ and $t=T$, the discrepancies $x(t)-\bar{x}$, $p_{x}(t)-\bar{p}_{x}$ and $u(t)-\bar{u}$ are bounded above by $\frac{1}{T}$, which is small as $T$ is large, but not exponentially small: it is weaker than (11), and we speak of a linear turnpike estimate. The second estimate in (22) gives a bound (uniform with respect to $T$ ) on the discrepancy $y(t)-\bar{y}(t)$ along the whole interval $[0, T]$ : it says that $y(t)$ remains at a uniform distance of $\bar{y}(t)=y^{0}+\operatorname{tg}(\bar{x}, \bar{u})$ as $t \in[0, T]$. Finally, the third estimate in (22) says that the constant $p_{y}-\bar{p}_{y}$ is (linearly) small as $T$ is large.

Of course, on specific examples it may happen that some of the components of the extremal triple $(x(\cdot), p(\cdot), u(\cdot))$ enjoy the exponential turnpike property, i.e., an estimate that is stronger than (22), without the term $\frac{1}{T}$ at the right-hand side. This is the case in the example (15): one can see in (17) that the component $x$ satisfies the linear turnpike estimate (22), while the components $p_{x}$ and $u$ satisfy the stronger exponential turnpike estimate (11).

Remark 4. As said earlier, to simplify exposition we have assumed that, for every $T \geqslant T_{0}$, the optimal control problem (12)-(13)-(14) has a unique solution, whose extremal lift is as well unique and normal, and that the static optimal control problem (21) has a unique solution, whose Lagrange multiplier is as well unique and normal. In case the optimal solution is not unique, or the extremal lift (still assumed to be normal) is not unique, the statement of Theorem (2) must be adapted as in [29]: we obtain in this case a result that is satisfied locally around the extremal triple.

Anyway, when the optimal control problem (12)-(13)-(14) is linear-quadratic,

$$
\begin{array}{lll}
\dot{x}(t)=\bar{A}_{1} x(t)+\bar{B}_{1} u(t), & x(0) \in M_{0}, & x(T) \in M \\
\dot{y}(t)=\bar{A}_{2} x(t)+\bar{B}_{2} u(t), & y(0)=y_{0}, & y(T)=y_{1}^{T}  \tag{23}\\
\min \int_{0}^{T}\left(x(t)^{\top} W x(t)+u(t)^{\top} U u(t)\right) d t & &
\end{array}
$$

i.e., when $f(x, u)=\bar{A}_{1} x+\bar{B}_{1} u, g(x, u)=\bar{A}_{2} x+\bar{B}_{2} u, f^{0}(x, u)=x^{\top} W x+u^{\top} U u$, the result of Theorem 2 is automatically global, by strict convexity, without having to assume any dissipativity property.

Theorem 2 is proved in Section 2. In Section 3 we give several optimal control problems that illustrate our main result and we provide numerical simulations showing evidence of the linear turnpike phenomenon.

## 2 Proof of Theorem 2

### 2.1 Preliminaries on the static optimal control problem

Let us consider the static optimal control problem (20), written in the form of the optimization problem (21): in the latter form, the optimization problem is "well-conditioned" in the sense that the constraints in $(x, u)$ are of the order of 1 as $T$ is large. This observation is important to derive the lemma below.

For the minimizer $(\bar{x}, \bar{u})$ (which is assumed to exist and to be unique), we denote by $\bar{p}_{x} \in \mathbb{R}^{n}$ and $\bar{p}_{y} \in \mathbb{R}^{p}$ the (normal) Lagrange multipliers respectively associated to the two constraints. Note that $\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)$ depends on $T$ but we do not add an index $T$ to keep better readability.

Lemma 2. We have

$$
\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)=\mathrm{O}(1)
$$

as $T \rightarrow+\infty$.

Proof. By the Lagrange multiplier rule, we have the optimality system

$$
F\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)=\left(\begin{array}{c}
d f(\bar{x}, \bar{u})^{\top} \bar{p}_{x}+d g(\bar{x}, \bar{u})^{\top} \bar{p}_{y}-d f^{0}(\bar{x}, \bar{u})^{\top} \\
f(\bar{x}, \bar{u}) \\
g(\bar{x}, \bar{u})
\end{array}\right)=\left(\begin{array}{c}
0 \\
0 \\
\frac{y_{1}^{T}-y_{0}}{T}
\end{array}\right)
$$

which is a system of $2 n+m+p$ equations with the $2 n+m+p$ unknowns $\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)$, where $F: \mathbb{R}^{2 n+m+p} \rightarrow \mathbb{R}^{2 n+m+p}$ is a $C^{1}$ mapping. The Jacobian of $F$ at the point $\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)$, which is the usual sensitivity matrix in optimization, is

$$
d F\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y_{1}}, \bar{p}_{y_{2}}\right)=\left(\begin{array}{cc|cc}
H_{x x} & H_{x u} & H_{x p_{x}} & H_{x p_{y}} \\
H_{u x} & H_{u u} & H_{u p_{x}} & H_{u p_{y}} \\
\hline H_{p_{x} x} & H_{p_{x} u} & 0 & 0 \\
H_{p_{y} x} & H_{p_{y} u} & 0 & 0
\end{array}\right)=\left(\begin{array}{c|c}
E_{1} & E_{2}^{\top} \\
\hline E_{2} & 0
\end{array}\right)
$$

Note that

$$
E_{2}=\left(\begin{array}{ll}
\frac{\partial f}{\partial x}(\bar{x}, \bar{u}) & \frac{\partial f}{\partial u}(\bar{x}, \bar{u}) \\
\frac{\partial g}{\partial x}(\bar{x}, \bar{u}) & \frac{\partial g}{\partial u}(\bar{x}, \bar{u})
\end{array}\right)=\left(\begin{array}{ll}
\bar{A} & B
\end{array}\right)
$$

and since $\operatorname{ker}\left(\bar{A}^{\top}\right) \cap \operatorname{ker}\left(\bar{B}^{\top}\right)=\{0\}$ by assumption, it follows immediately that the matrix $E_{2}$ is surjective.

Besides, let us prove that the matrix $E_{1}$ is invertible. If $E_{1}\binom{\xi_{1}}{\xi_{2}}=0$ then $H_{x x} \xi_{1}+H_{x u} \xi_{2}=0$ and $H_{u x} \xi_{1}+H_{u u} \xi_{2}=0$. From the latter equation, we get $\xi_{2}=-H_{u u}^{-1} H_{u x} \xi_{1}$, and plugging into the first equation yields $\left(H_{x x}-H_{x u} H_{u u}^{-1} H_{u x}\right) \xi_{1}=0$. Since, by assumption, the matrix $W=-H_{x x}+H_{x u} H_{u u}^{-} 1 H_{u x}$ is symmetric positive definite, we infer that $\xi_{1}=0$, and then $\xi_{2}=0$. Therefore $E_{1}$ is invertible.

Now, since $E_{1}$ is invertible and $E_{2}$ is surjective, we conclude that $d F\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)$ is invertible. The assumption on $y_{1}^{T}$ done in the statement of Theorem 2 states that $\frac{y_{1}^{T}-y_{0}}{T}=\mathrm{O}(1)$ as $T \rightarrow+\infty$. Hence, we have to solve the nonlinear system of equations $F\left(\bar{x}, \bar{u}, \bar{p}_{x}, \bar{p}_{y}\right)=\mathrm{O}(1)$. The lemma follows by applying the implicit function theorem.

### 2.2 Perturbation variables

Let us now introduce perturbation variables, by setting
$\delta x(t)=x(t)-\bar{x}, \quad \delta y(t)=y(t)-\bar{y}(t), \quad \delta p_{x}(t)=p_{x}(t)-\bar{p}_{x}, \quad \delta p_{y}=p_{y}-\bar{p}_{y}, \quad \delta u(t)=u(t)-\bar{u}$.
In these variables, we have the terminal conditions

$$
\begin{equation*}
\delta x(0) \in M_{0}-\bar{x}, \quad \delta y(0)=0, \quad \delta x(T) \in M_{1}-\bar{x}, \quad \delta y(T)=0 \tag{24}
\end{equation*}
$$

Using Lemma 2, we note the important fact that all terminal values $\delta x(0), \delta y(0), \delta p_{x}(0), \delta p_{y}$, $\delta x(T), \delta y(T), \delta p_{x}(T)$, are at most of the order of 1 as $T$ is large.

Linearizing (19) and then (18) along $\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y}, \bar{u}\right)$, we get, at the first order (i.e., neglecting the terms of order $\geqslant 2$ ),

$$
\delta u(t)=U^{-1}\left(H_{u x} \delta x(t)+\bar{B}_{1}^{\top} \delta p_{x}(t)+\bar{B}_{2}^{\top} \delta p_{y}\right)
$$

and then, still at the first order,

$$
\begin{aligned}
\delta \dot{x}(t) & =A_{1} \delta x(t)+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top} \delta p_{x}(t)+\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top} \delta p_{y} \\
\delta \dot{y}(t) & =A_{2} \delta x(t)+\bar{B}_{2} U^{-1} \bar{B}_{1}^{\top} \delta p_{x}(t)+\bar{B}_{2} U^{-1} \bar{B}_{2}^{\top} \delta p_{y} \\
\delta \dot{p}_{x}(t) & =W \delta x(t)-A_{1}^{\top} \delta p_{x}(t)-A_{2}^{\top} \delta p_{y}
\end{aligned}
$$

where we recall that $\delta p_{y}$ is constant. Setting

$$
z(t)=\binom{\delta x(t)}{\delta p_{x}(t)}
$$

and defining the matrices

$$
M=\left(\begin{array}{cc}
A_{1} & \bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}  \tag{25}\\
W & -A_{1}^{\top}
\end{array}\right), \quad V=\binom{-\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}}{A_{2}^{\top}}, \quad L=\left(\begin{array}{ll}
A_{2} & \bar{B}_{2} U^{-1} \bar{B}_{1}^{\top}
\end{array}\right)
$$

we write the above system in the form

$$
\begin{align*}
\dot{z}(t) & =M z(t)-V \delta p_{y}  \tag{26}\\
\delta \dot{y}(t) & =L z(t)+\bar{B}_{2} U^{-1} \bar{B}_{2}^{\top} \delta p_{y} \tag{27}
\end{align*}
$$

When $H_{u x}=0$, this corresponds exactly to the extremal system associated with the linearquadratic optimal control problem (23).

We have the following spectral property for the matrix $M$.
Lemma 3. The matrix $M$ is hyperbolic, i.e., all (complex) eigenvalues of $M$ have a nonzero real part.
Proof. Let $\mu \in \mathbb{R}$ and let $\xi \in \mathbb{C}^{2 n}$ be such that $(M-i \mu I) \xi=0$. Let us prove that $\xi=0$. Denoting by $\left(\xi_{1}, \xi_{2}\right)$ the coordinates of $z$, we have

$$
\begin{align*}
& \left(A_{1}-i \mu I_{n}\right) \xi_{1}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top} \xi_{2}=0  \tag{28}\\
& W \xi_{1}-\left(A_{1}^{\top}+i \mu I_{n}\right) \xi_{2}=0
\end{align*}
$$

We stress that all matrices above are real, but the coordinates $\left(\xi_{1}, \xi_{2}\right)$ are complex. The second equation in (28) gives $\xi_{1}=W^{-1}\left(A_{1}^{\top}+i \mu I_{n}\right) \xi_{2}$, and plugging this expression in the first equation of (28) yields

$$
\left(A_{1}-i \mu I_{n}\right) W^{-1}\left(A_{1}^{\top}+i \mu I_{n}\right) \xi_{2}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top} \xi_{2}=0
$$

and multiplying this equality, to the left, by the complex conjugate $\bar{\xi}_{1}^{\top}$, gives

$$
\begin{equation*}
\left\|W^{-1 / 2}\left(A_{1}^{\top}+i \mu I_{n}\right) \xi_{2}\right\|^{2}+\left\|U^{-1 / 2} \bar{B}_{1}^{\top} \xi_{2}\right\|^{2}=0 \tag{29}
\end{equation*}
$$

Note that here, we have used that the complex conjugate of $A_{1}-i \mu I_{n}$ is ${\overline{\left(A_{1}-i \mu I_{n}\right)}}^{\top}=A_{1}^{\top}+i \mu I_{n}$. The norms and scalar product used in (29) are Hermitian. Now, (29) implies that $\left(A_{1}^{\top}+i \mu I_{n}\right) \xi_{2}=0$ and $\bar{B}_{1}^{\top} \xi_{2}=0$, i.e.,

$$
\begin{equation*}
\xi_{2} \in \operatorname{ker}\left(A_{1}^{\top}+i \mu I_{n}\right) \cap \operatorname{ker}\left(\bar{B}_{1}^{\top}\right) \tag{30}
\end{equation*}
$$

Since the pair $\left(A_{1}, \bar{B}_{1}\right)$ satisfies the Kalman condition (see Foonote 1 in Theorem 1), we have $\operatorname{ker}\left(A_{1}^{\top}-\lambda I_{n}\right) \cap \operatorname{ker}\left(\bar{B}_{1}^{\top}\right)=\{0\}$ for every $\lambda \in \mathbb{C}$ (Hautus test, see, e.g., [23, 25]). Then, (30) gives $\xi_{2}=0$.

The second equation of (28) gives then $\xi_{1}=0$. We conclude that $\xi=0$.
Since $M$ is hyperbolic, first of all, it is invertible and then (26) can be written as

$$
\frac{d}{d t}\left(z(t)-M^{-1} V \delta p_{y}\right)=M\left(z(t)-M^{-1} V \delta p_{y}\right)
$$

Second, there exists a real-valued invertible matrix $P$ of size $2(n+p)$ such that

$$
P^{-1} M P=\left(\begin{array}{cc}
M_{1} & 0 \\
0 & M_{2}
\end{array}\right)
$$

where all (complex) eigenvalues of the real-valued matrix $M_{1}$ have a negative real part and all (complex) eigenvalues of the real-valued matrix $M_{2}$ have a positive real part. Setting ${ }^{3}$

$$
z(t)-M^{-1} V \delta p_{y}=P\binom{v_{1}(t)}{v_{2}(t)}
$$

we have then $\dot{v}_{1}(t)=M_{1} v_{1}(t)$ and $\dot{v}_{2}(t)=M_{2} v_{2}(t)$. Since all complex eigenvalues of $M_{1}$ (resp., of $M_{2}$ ) have a negative (resp., positive) real part, there exist $C>0$ and $\nu>0($ not depending on $T)$ such that $\left\|v_{1}(t)\right\| \leqslant C e^{-\nu t}$ and $\left\|v_{2}(t)\right\| \leqslant C e^{-\nu(T-t)}$ for every $t \geqslant 0$.

Noting (by using (24)) that $z(0)-M^{-1} V \delta p_{y}=\mathrm{O}(1)$ and $z(T)-M^{-1} V \delta p_{y}=\mathrm{O}(1)$ as $T \rightarrow+\infty$, we get (taking a larger constant $C$ if necessary) that

$$
\begin{equation*}
\left\|z(t)-M^{-1} V \delta p_{y}\right\| \leqslant C\left(e^{-\nu t}+e^{-\nu(T-t)}\right) \quad \forall t \in[0, T] \tag{31}
\end{equation*}
$$

Now, integrating (27), using that $\int_{0}^{T} e^{-\nu t} d t \leqslant \frac{1}{\nu}$ for every $T>0$, and defining the square matrix of size $p$

$$
\begin{equation*}
R=L M^{-1} V+\bar{B}_{2} U^{-1} \bar{B}_{2}^{\top} \tag{32}
\end{equation*}
$$

we infer from (31) that (taking a larger constant $C$ if necessary)

$$
\begin{equation*}
\left\|\delta y(t)-t R \delta p_{y}\right\| \leqslant C \quad \forall t \in[0, T] \tag{33}
\end{equation*}
$$

Since we must have $\delta y(T)=0$ (this comes from (24), itself coming from the final condition $y(T)=y_{1}^{T}$ on the component $y$ ), we infer that

$$
\begin{equation*}
\left\|R \delta p_{y}\right\| \leqslant \frac{C}{T} \tag{34}
\end{equation*}
$$

We have the following instrumental fact.
Lemma 4. The matrix $R$ defined by (32) is symmetric positive definite, and thus is invertible.
Lemma 4 is proved in Section 2.3.
It follows from (34) and from Lemma 4 that (taking a larger constant $C$ if necessary) $\left\|\delta p_{y}\right\| \leqslant \frac{C}{T}$. Let us then finish the proof of Theorem 2. Since $\delta p_{y_{2}}=\mathrm{O}\left(\frac{1}{T}\right)$, taking a larger constant $C$ if necessary, we infer from (31) and (33) that

$$
\|z(t)\| \leqslant C\left(\frac{1}{T}+e^{-\nu t}+e^{-\nu(T-t)}\right), \quad\|\delta y(t)\| \leqslant C \quad \forall t \in[0, T]
$$

which gives the estimates stated in Theorem 2 at the first order around $\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y}\right)$.
Until now we have kept only the first-order terms, in the linearization process. To take into account the higher-order terms, as well as the general terminal conditions $x(0) \in M_{0}, x(T) \in M_{1}$, $y(0)=y_{0}, y(T)=y_{1}$ and the corresponding transversality conditions (4), we have to prove that, locally around the minimizer of the static problem, the corresponding shooting problem is locally well posed. Since this part of the proof is exactly the same as in [29, end of Section 3.2], we do not reproduce it here.

[^3]From local to global. At this step, the turnpike property has been established locally around $\left(\bar{x}, \bar{y}(t), \bar{p}_{x}, \bar{p}_{y}\right)$. Let us now prove that the turnpike estimate is global. We sketch hereafter the argument. It suffices to prove that the unique solution $\left(x(\cdot), y(\cdot), p_{x}(\cdot), p_{y}(\cdot), u(\cdot)\right)$ of the optimal control problem (12)-(13)-(14) passes arbitrarily close to the turnpike set $\left(\bar{x}, \bar{y}(\cdot), \bar{p}_{x}, \bar{p}_{y}\right)$ as $T$ is large: indeed, along such a time interval, the turnpike inequality is true thanks to the above proof; and thus, taking larger constants if necessary (which is possible by compactness, see Assumption (i) in Theorem 1), the turnpike inequality remains true along the whole interval $[0, T]$, uniformly with respect to $T$. By uniqueness of the extremal lift (see Assumptions (i) and (ii)), it suffices to show that $x(\cdot)$ passes arbitrarily close to $\bar{x}$ as $T$ is large. This property follows from the strict dissipativity assumption (iii) and from the controllability assumption (iv). Indeed, the strict dissipativity inequality (10) implies that

$$
f^{0}(\bar{x}, \bar{u}) \leqslant \frac{1}{T} \int_{0}^{T} f^{0}(x(t), u(t)) d t+\frac{S(x(0))-S(x(T))}{T}-\frac{1}{T} \int_{0}^{T} \alpha(\|x(t)-\bar{x}\|,\|u(t)-\bar{u}\|) d t
$$

Let us prove by contradiction that $x(\cdot)$ passes near $\bar{x}$. Otherwise, this means that the term $\frac{1}{T} \int_{0}^{T} \alpha(\|x(t)-\bar{x}\|,\|u(t)-\bar{u}\|) d t$ is of the order of 1 and we infer from the above inequality that $f^{0}(\bar{x}, \bar{u}) \leqslant \frac{1}{T} \int_{0}^{T} f^{0}(x(t), u(t)) d t-\eta$ for some $\eta>0$. Besides, by the controllability assumption (iv), there exists a trajectory $x_{1}(\cdot)$, solution of (1) on $[0, T]$, such that $x_{1}(0) \in M_{0}, x_{1}\left(\delta_{1}\right)=\bar{x}$, then $x_{1}(t) \equiv \bar{x}$ on $\left[\delta_{1}, T-\delta_{2}\right]$, and $x_{1}(T) \in M_{1}$. The cost of this trajectory is

$$
\begin{aligned}
\frac{1}{T} \int_{0}^{T} f^{0}\left(x_{1}(t), u_{1}(t)\right) d t=\frac{1}{T} \int_{0}^{\delta_{1}} f^{0}\left(x_{1}(t), u_{1}(t)\right) d t & +\frac{1}{T} \int_{T-\delta_{2}}^{T} f^{0}\left(x_{1}(t), u_{1}(t)\right) d t \\
& +\frac{T-\delta_{1}-\delta_{2}}{T} f^{0}(\bar{x}, \bar{u})=f^{0}(\bar{x}, \bar{u})+\mathrm{O}\left(\frac{1}{T}\right)
\end{aligned}
$$

as $T$ is large, and hence we have obtained a trajectory $x_{1}(\cdot)$ whose cost is (strictly) less than the cost of the optimal trajectory $x(\cdot)$ : this is a contradiction.

Remark 5. In (12)-(13)-(14), the dynamics in $y$ is written as $\dot{y}(t)=h(x(t), y(t), u(t))$ with $h(x, y, u)=g(x, u)$ : the mapping $h$ does not depend on $y$. Our strategy of proof would not apply to a general mapping $h$, depending on $y$. Indeed, otherwise, when linearizing $\dot{y}(t)=h(x(t), y(t), u(t))$ along the path $(\bar{x}, \bar{y}(t), \bar{u})$, we would get $\delta \dot{y}(t)=A_{2} \delta x(t)+\bar{B}_{2} \delta u(t)+Q(t) \delta y(t)$ where $Q(t)=$ $\frac{\partial h}{\partial y}(\bar{x}, \bar{y}(t), \bar{u})$ depends on $t$ and thus the matrix $M$ defined by (25) would as well depend on $t$ (anyway, only through the part in $Q$ ). The argument of hyperbolicity, then, cannot be applied in general, in particular we think that it may fail whenever the eigenvectors of $Q(t)$ are oscillating too fast. Nevertheless, we think that the technique of proof may be extended to the case where $t \mapsto Q(t)$ is slowly-varying, as in the context of control or stabilization by quasi-static deformation (see [8] and references therein). We leave this issue open.

### 2.3 Proof of Lemma 4

Let us first express the matrix $R$, defined by (32), in a more explicit way. Recall that

$$
R=\left(A_{2} \bar{B}_{2} U^{-1} \bar{B}_{1}^{\top}\right) M^{-1}\binom{-\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}}{A_{2}^{\top}}+\bar{B}_{2} U^{-1} \bar{B}_{2}^{\top} \quad \text { with } \quad M=\left(\begin{array}{cc}
A_{1} & \bar{B}_{1} U^{-1} \bar{B}_{1}^{\top} \\
W & -A_{1}^{\top}
\end{array}\right)
$$

We set

$$
\binom{C_{1}}{C_{2}}=M^{-1}\binom{-\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}}{A_{2}^{\top}}
$$

so that

$$
\begin{align*}
A_{1} C_{1}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top} C_{2} & =-\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top} \\
W C_{1}-A_{1}^{\top} C_{2} & =A_{2}^{\top} \tag{35}
\end{align*}
$$

Since $W$ is invertible, we infer from the second equation of (35) that $C_{1}=W^{-1} A_{1}^{\top} C_{2}+W^{-1} A_{2}^{\top}$, and plugging this expression into the first equation of (35) gives

$$
\begin{equation*}
\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right) C_{2}=-A_{1} W^{-1} A_{2}^{\top}-\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top} \tag{36}
\end{equation*}
$$

Lemma 5. The matrix $A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}$ is invertible.
Proof. Let $\xi \in \mathbb{R}^{n}$ be such that $\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right) \xi=0$. Multiplying to the left by $\xi^{\top}$, we obtain $\left\|W^{-1 / 2} A_{1}^{\top} \xi\right\|^{2}+\left\|U^{-1 / 2} \bar{B}_{1}^{\top} \xi\right\|^{2}=0$ and hence $\xi \in \operatorname{ker}\left(A_{1}^{\top}\right) \cap \operatorname{ker}\left(\bar{B}_{1}^{\top}\right)$. As in the proof of Lemma 3, since the pair $\left(A_{1}, \bar{B}_{1}\right)$ satisfies the Kalman condition, the latter intersection is $\{0\}$. The claim follows.

We infer from (36) and from Lemma 5 that

$$
C_{2}=-\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right)^{-1}\left(A_{1} W^{-1} A_{2}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}\right)
$$

and thus

$$
C_{1}=-W^{-1} A_{1}^{\top}\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right)^{-1}\left(A_{1} W^{-1} A_{2}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}\right)+W^{-1} A_{2}^{\top}
$$

We conclude that

$$
\begin{aligned}
& R=-\left(A_{1} W^{-1} A_{2}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}\right)^{\top}\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right)^{-1}\left(A_{1} W^{-1} A_{2}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{2}^{\top}\right) \\
&+A_{2} W^{-1} A_{2}^{\top}+\bar{B}_{2} U^{-1} \bar{B}_{2}^{\top}
\end{aligned}
$$

(note that $R$ is symmetric) that we can write as the sum of two symmetric matrices

$$
\begin{aligned}
R= & A_{2} W^{-1 / 2}\left(-W^{-1 / 2} A_{1}^{\top}\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right)^{-1} A_{1} W^{-1 / 2}+I_{n}\right) W^{-1 / 2} A_{2}^{\top} \\
& +\bar{B}_{2} U^{-1 / 2}\left(-U^{-1 / 2} \bar{B}_{1}^{\top}\left(A_{1} W^{-1} A_{1}^{\top}+\bar{B}_{1} U^{-1} \bar{B}_{1}^{\top}\right)^{-1} \bar{B}_{1} U^{-1 / 2}+I_{m}\right) U^{-1 / 2} \bar{B}_{2}^{\top} \\
= & \tilde{A}_{2}\left(-\tilde{A}_{1}^{\top}\left(\tilde{A}_{1} \tilde{A}_{1}^{\top}+\tilde{B}_{1} \tilde{B}_{1}^{\top}\right)^{-1} \tilde{A}_{1}+I_{n}\right) \tilde{A}_{2}^{\top}+\tilde{B}_{2}\left(-\tilde{B}_{1}^{\top}\left(\tilde{A}_{1} \tilde{A}_{1}^{\top}+\tilde{B}_{1} \tilde{B}_{1}^{\top}\right)^{-1} \tilde{B}_{1}+I_{m}\right) \tilde{B}_{2}^{\top}
\end{aligned}
$$

where we have set $\tilde{A}_{i}=A_{i} W^{-1 / 2}$ and $\tilde{B}_{i}=\bar{B}_{i} U^{-1 / 2}$ for $i=1,2$, and we are going to prove hereafter that the above two matrices are symmetric positive definite, so that $R$ is itself symmetric positive definite.

We have the following general result.
Lemma 6. Let $A$ be an arbitrary real-valued square matrix of size $n$, and let $B$ be an arbitrary realvalued matrix of size $n \times m$, where $n$ and $m$ are arbitrary nonzero integers. If $\operatorname{ker}\left(A^{\top}\right) \cap \operatorname{ker}\left(B^{\top}\right)=$ $\{0\}$ then

$$
\begin{equation*}
B^{\top}\left(A A^{\top}+B B^{\top}\right)^{-1} B \preceq I_{m} \tag{37}
\end{equation*}
$$

meaning that the symmetric matrix $I_{m}-B^{\top}\left(A A^{\top}+B B^{\top}\right)^{-1} B$ is positive semi-definite.

Proof. Since $\operatorname{ker}\left(A^{\top}\right) \cap \operatorname{ker}\left(B^{\top}\right)=\{0\}$, the matrix $A A^{\top}+B B^{\top}$ is invertible (same argument as in Lemma 5).

When $B$ is invertible, the result is obvious: starting from

$$
A A^{\top}+B B^{\top} \succeq B B^{\top}
$$

taking the inverse and multiplying to the left by $B^{\top}$ and to the right by $B$, we obtain (37).
When $B$ is not invertible, we follow anyway the above reasoning, adding $\varepsilon I_{m}$ for $\varepsilon>0$ to recover an appropriate invertibility property: starting from

$$
A A^{\top}+B B^{\top}+\varepsilon I_{m} \succeq B B^{\top}+\varepsilon I_{m}
$$

taking the inverse and multiplying to the left by $B^{\top}$ and to the right by $B$, we obtain

$$
B^{\top}\left(A A^{\top}+B B^{\top}+\varepsilon I_{m}\right)^{-1} B \preceq B^{\top}\left(B B^{\top}+\varepsilon I_{m}\right)^{-1} B
$$

Now, we use the general fact that

$$
B^{\top}\left(B B^{\top}+\varepsilon I_{m}\right)^{-1} B \underset{\varepsilon \rightarrow 0}{\longrightarrow} I_{m}
$$

which is a consequence of the Tikhonov regularization in the Moore-Penrose pseudo-inverse theory (see [24]). Then, taking the limit $\varepsilon \rightarrow 0$ gives (37).

Noting that $\operatorname{ker}\left(\tilde{A}_{1}^{\top}\right) \cap \operatorname{ker}\left(\tilde{B}_{1}^{\top}\right)=\{0\}$ (this comes again from the Kalman condition on the pair $\left(A_{1}, \bar{B}_{1}\right)$ ), we infer from Lemma 6 that

$$
\tilde{A}_{1}^{\top}\left(\tilde{A}_{1} \tilde{A}_{1}^{\top}+\tilde{B}_{1} \tilde{B}_{1}^{\top}\right)^{-1} \tilde{A}_{1} \preceq I_{n} \quad \text { and } \quad \tilde{B}_{1}^{\top}\left(\tilde{A}_{1} \tilde{A}_{1}^{\top}+\tilde{B}_{1} \tilde{B}_{1}^{\top}\right)^{-1} \tilde{B}_{1} \preceq I_{m}
$$

Therefore $R \succeq 0$, i.e., $R$ is positive semi-definite.

## 3 Examples

In this section, we illustrate the linear turnpike phenomenon on several practical examples, with numerical simulations. Hereafter, we compute numerically the optimal solutions either with a direct method, using AMPL (see [12]) combined with IpOpt (see [30]), following a full discretization of the problem, or with the variant of the shooting method introduced in [29]. In both cases, we initialize the numerical method at the turnpike, thus guaranteeing its convergence.

### 3.1 The Zermelo problem

The Zermelo problem is a famous optimal control problem, often used as a simple model or exercise to illustrate theory or practice of optimal control (see [25]). In this problem, we consider a boat moving with constant speed along a river of constant width $\ell$, in which there is a current $c(y)$. The movement of the center of mass of the boat is governed by the control system

$$
\begin{array}{llrl}
\dot{x}(t) & =v \cos u(t)+c(y(t)), & x(0) & =0, \\
\dot{y}(t) & =v \sin u(t), & y(0) & =0,
\end{array}
$$

where $v>0$ is the constant speed, and the control is the angle $u(t)$ of the axis of the boat with respect to the axis $(0 x)$ (see Figure 2). Given some $T>0$ fixed, we consider the optimal


Figure 1: Zermelo problem.
control problem of steering in time $T$ this control system from $M_{0}=\{(0,0)\}$ to the final set $M_{1}=\{y=\ell\}$ (upper river), by maximizing the lateral deport $x(T)$, which is equivalent to minimizing $-\int_{0}^{T}(v \cos u(t)+c(y(t))) d t$.

We take $T=20, v=1, L=2$ and $c(y)=3+0.2 y(L-y)$. The current in the river is so strong that we always have $\dot{x}(t)>0$, whatever the control may be; hence $x(t)$ is bound to be increasing.


Figure 2: Numerical results for the Zermelo problem
Solving the "static" problem (20) yields the turnpike trajectory $\bar{x}(t)=5 t, \bar{y}=1, \bar{u}=0$. We let the reader check that the assumptions of Theorem 2 are satisfied. We do not give any details on the application of the Pontryagin maximum principle to this problem, which is easy (see, e.g., [25]). We just mention that $p_{x}=1$ is constant. Here, our objective is to show evidence of the
linear turnpike phenomenon, by assuming that $T$ is large. The numerical results are reported on Figure 2.

### 3.2 Optimal control model of a runner

Consider the runner model developed in [1, 2, 3], originating from [18]:

$$
\begin{align*}
\dot{x}(t) & =v(t) & x(0)=0, \quad x\left(t_{f}\right) \\
\dot{v}(t) & =-\frac{v(t)}{\tau}+f(t) & v(0)=v^{0}  \tag{38}\\
\dot{e}(t) & =\sigma-f(t) v(t) & e(0)=e^{0}, \quad e\left(t_{f}\right) \\
\dot{f}(t) & =\gamma\left(u(t)\left(F_{\max }-f(t)\right)-f(t)\right) & 0 \leqslant f(t) \leqslant F_{\max }
\end{align*}
$$

where $d>0$ is the prescribed distance to run, $e^{0}>0$ is the initial energy, $\tau>0$ is the friction coefficient related to the runner's economy, $\sigma>0$ is a constant standing for the energetic equivalent of the oxygen uptake $V O_{2}, F_{\max }>0$ is a threshold upper bound for the force $f(t), \gamma>0$ is the time constant of motor activation and $u(t) \in[-M, M]$ is the neural drive which is the control, where $M>0$ is some control bound. Here, $x(t)$ is the distance travelled at time $t$ by the runner, $v(t)$ the velocity, $e(t)$ the anaerobic energy, and $f(t)$ the propulsive force per unit mass. Actually, the typical values of the initial energy $e_{0}$ are such that the energy variable $e(t)$ is decreasing. The minimization criterion is

$$
\begin{equation*}
\min \left(t_{f}+\frac{\alpha}{2} \int_{0}^{t_{f}} u(t)^{2} d t\right) \tag{39}
\end{equation*}
$$

where $\alpha>0$ is a fixed constant and the final time $t_{f}$ is free.
Although, in the optimal control problem (38)-(39), $t_{f}, v\left(t_{f}\right), f(0)$ and $f\left(t_{f}\right)$ are let free, the problem can be reparametrized by the distance $s$, with the change of variable $t^{\prime}(s)=\frac{d t}{d s}=\frac{1}{v(s)}$. In terms of $s$, the optimal control problem is

$$
\begin{array}{ll}
v^{\prime}(s)=-\frac{1}{\tau}+\frac{f(s)}{v(s)} & v(0)=v^{0} \\
e^{\prime}(s)=\frac{\sigma}{v(s)}-f(s) & e(0)=e^{0}, \quad e(d)= \\
f^{\prime}(s)=\frac{\gamma}{v(s)}\left(u(s)\left(F_{\max }-f(s)\right)-f(s)\right) & 0 \leqslant f(s) \leqslant F_{\max } \\
\min \int_{0}^{d} \frac{1}{v(s)}\left(1+\frac{\alpha}{2} u(s)^{2}\right) d s &
\end{array}
$$

and now fits in the general framework developed in this paper.
Solving the "static" problem leads to $\bar{e}(t)=e_{0}+(\sigma-\bar{f} \bar{v}) t$ (we re-express it in function of $t$ ) where

$$
\bar{v}=\frac{\tau}{2 d}\left(e_{0}+\sqrt{e_{0}^{2}+4 \frac{\sigma d^{2}}{\tau}}\right), \quad \bar{f}=\frac{\bar{v}}{\tau}=\frac{1}{2 d}\left(e_{0}+\sqrt{e_{0}^{2}+4 \frac{\sigma d^{2}}{\tau}}\right), \quad \bar{u}=\frac{\bar{f}}{F_{\max }-\bar{f}} .
$$

For the numerical simulations, we take $d=1500, \tau=0.932, \sigma=22, \alpha=10^{-5}, F_{\max }=8$, $\gamma=0.0025, v_{0}=3, e_{0}=4651$. The distance run ( 1500 meters) is large enough so that we observe the linear turnpike phenomenon. The numerical results are reported on Figure 3.

This runner optimal control problem has actually been the initial point for the present paper, and the author warmly thanks Amandine Aftalion for having raised such an interesting problem.


Figure 3: Numerical results for the Zermelo problem

## A Appendix

This appendix is devoted to illustrate the comments done in Remark 1. We take a very simple example and we give numerical simulations showing the competition between two global turnpikes, or, at the level of the initialization of numerical methods, between local and global turnpikes.

We consider the one-dimensional optimal control problem

$$
\begin{align*}
& \dot{x}=-3 x+3 x^{3}+u, \quad x(0)=x_{0}, \quad x(T)=x_{f} \text { or free } \\
& \min \int_{0}^{T}\left(\left(x(t)-x_{d}\right)^{2}+\left(u(t)-u_{d}\right)^{2}\right) d t \tag{40}
\end{align*}
$$

The corresponding static problem is

$$
\begin{equation*}
\min _{(x, u) \mid-3 x+3 x^{3}+u=0}\left(\left(x-x_{d}\right)^{2}+\left(u-u_{d}\right)^{2}\right) . \tag{41}
\end{equation*}
$$

## A. 1 Competition between two (global) turnpikes

We take $x_{d}=1$ and $u_{d}=3.47197$. The choice of $u_{d}$ is done so that the static problem (41) has two global minima, at $\bar{x}_{1}=-1.347372066$ and $\bar{x}_{2}=0.5939615956$, see Figure 4 .


Figure 4: Plot of the function $x \mapsto(x-1)^{2}+\left(3 x-3 x^{3}-3.47197\right)^{2}$
On Figure 5, in dashed blue, we have computed the optimal trajectory with $x_{0}=-5, x_{f}=-1$, $T=10$ : we observe a turnpike phenomenon around $\bar{x}_{1}$. In solid red, we have computed the optimal trajectory with $x_{0}=2, x_{f}=1, T=10$ : we observe a turnpike phenomenon around $\bar{x}_{2}$.


Figure 5: Global turnpikes around $\bar{x}_{1}$ and $\bar{x}_{2}$.
The fact that the turnpike phenomenon is either around $\bar{x}_{1}$ or around $\bar{x}_{2}$ depends on the terminal conditions. For instance, if $x_{0}$ and $x_{f}$ are close to $\bar{x}_{1}$ (resp., $\bar{x}_{2}$ ) then the optimal trajectory will make a turnpike around $\bar{x}_{1}$ (resp., $\bar{x}_{2}$ ). But when the terminal conditions are farther, it is not clear to predict the behavior of the optimal trajectory.

Since the minima are global, we expect to observe a competition between both turnpikes, depending on the terminal conditions (see [22]). Let us provide some numerical simulations illustrating this competition. To facilitate the understanding, we consider the problem with $x_{f}$ free. It is interesting to note that the numerical result strongly depends on the initialization of the numerical method. Here, to compute numerically the optimal solutions of (40), we use AMPL (see [12]) combined with IpOpt (see [30]): the trajectory and the controls are discretized (the control is piecewise constant and the trajectory is piecewise linear, on a given subdivision that is chosen fine enough), and we initialize the trajectory with the same constant value over all the subdivision.

Then according to the value of this initialization, we can make emerge such or such turnpike, and all solutions are anyway optimal.


Figure 6: $x_{0}=-2, x_{f}$ free, $T=20$.

On Figure 6, we have taken $T=20, x_{0}=-2$ (and $x_{f}$ is let free). In dashed blue, we have initialized the trajectory to the constant trajectory $\bar{x}_{1}$, and we then obtain an optimal trajectory which stays essentially near $\bar{x}_{1}=-1.347372066$, with a kind of "hesitation" towards $\bar{x}_{2}=0.5939615956$ near the end. In solid red, we have initialized the trajectory to the constant trajectory $\bar{x}_{2}$, and we obtain a trajectory staying essentially near $\bar{x}_{2}$, with a kind of "hesitation" towards $\bar{x}_{1}$ near the beginning. We stress that the two solutions are optimal: both have a cost $C \simeq 6.2822$. This is so, because we have taken $T=20$, which is large enough. We could make emerge other similar trajectories, which "hesitate" between the two turnpikes. All of them are optimal, or, at least, "quasi-optimal" (there is a small error due to switches from one turnpike to the other).

## A. 2 Local versus global turnpike

We take $x_{d}=1$ and $u_{d}=1$. The choice of $u_{d}$ is now such that the static problem (41) has a unique global solution $\bar{x}=0.781538640850898$, see Figure 7. But it has also a local solution $\bar{x}_{l o c}=-1.10551208794920$.

The global minimum $\bar{x}$ is a global turnpike, while the local minimum $\bar{x}$ is a local turnpike. In the numerical simulations, when performing either an optimization or a Newton method (shooting), we compute local solutions. Hence, we must expect that the numerical results depend on the initialization. To check global optimality, we have to compare the costs.

On Figure 8, in dashed blue, we have initialized the code with the constant trajectory $\bar{x}_{l o c}$. We observe a turnpike around $\bar{x}_{l o c}=-1.10551208794920$. The cost is $C \simeq 7.008$. But this is a local turnpike only. This trajectory that we obtain is only locally optimal, and is not globally optimal. In solid red, we have initialized the code with the constant trajectory $\bar{x}$. We observe a turnpike around $\bar{x}=0.781538640850898$. The cost is $C \simeq 3.825$ and is lower than the one of the previous one. Here, we have actually computed the globally optimal trajectory. This is the global turnpike.

It is also interesting to see what happens if we take $T$ much smaller. Let us take $T=2$. On Figure 9, in dashed blue, we have initialized the code with the constant trajectory $\bar{x}$. We observe a trend to the turnpike around $\bar{x}=0.781538640850898$. The cost is $C \simeq 17.792$. But this trajectory,


Figure 7: Plot of the function $x \mapsto(x-1)^{2}+\left(3 x-3 x^{3}-1\right)^{2}$


Figure 8: $x_{0}=-2, x_{f}=-1, T=10$.
now, is not globally optimal. In solid red, we have initialized the code with the constant trajectory $\bar{x}_{l o c}$. We observe a turnpike around $\bar{x}_{l o c}=-1.347372066$. The cost is $C \simeq 16.322$. It is the globally optimal trajectory.

We can search a time $2<T_{c}<10$ for which both previous initializations give equivalent turnpikes around $\bar{x}_{l o c}$ and $\bar{x}$ (with same cost). This is done hereafter. What is important is that, in large time $T$, we have indeed the global turnpike around $\bar{x}$.

On Figure 10, at the left, the code is initialized with $x(t) \equiv-1.1$, in order to promote the turnpike around $\bar{x}_{l o c}=-1.347372066$. At the right, the code is initialized with $x(t) \equiv 0.78$, in order to promote the turnpike around $\bar{x}=0.781538640850898$. We have computed trajectories for the following successives values of $T: 2.5,2.7,2.9,3.1,3.3$. We observe that the blue and cyan trajectories at the top are optimal (see the value of their cost); and that the red, green and black trajectories at the bottom are optimal. The bifurcation occurs around $T=2.9$.

Finally, on Figure 11, we represent the global optimal trajectory. For $T \lesssim 2.9$, the global


Figure 9: $x_{0}=-2, x_{f}=-1, T=2$.


Figure 10: $x_{0}=-2, x_{f}=-1$
optimal trajectory makes a turnpike around $\bar{x}_{l o c}=-1.347372066$, which is a local minimizer of the optimal static problem. For $T \geqslant 2.9$ we have a bifurcation and the global optimal trajectory makes a turnpike around $\bar{x}=0.781538640850898$, which is the global minimizer of the optimal static problem.

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Figure 11: $x_{0}=-2, x_{f}=-1$
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[^1]:    ${ }^{1}$ Equivalently, the pair $\left(A_{1}, \bar{B}_{1}\right)$ satisfies the Kalman condition.

[^2]:    ${ }^{2}$ This condition is satisfied if the pair $(\bar{A}, B)$ satisfies the Kalman condition, i.e., $\operatorname{rank}\left(B, \bar{A} B, \ldots, \bar{A}^{n+p-1} B\right)=$ $n+p$, or in other words, if the linearized system at $(\bar{x}, \bar{u})$ is controllable.

[^3]:    ${ }^{3}$ Actually, since the matrix $M$ is Hamiltonian (see [29] for some remarks on this issue), the size of $v_{1}$ is equal to the size of $v_{2}$ (equal to $n+p$ ).

