# Boundary stabilization in finite time of one-dimensional linear hyperbolic balance laws with coefficients depending on time and space 

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#### Abstract

In this article we are interested in the boundary stabilization in finite time of one-dimensional linear hyperbolic balance laws with coefficients depending on time and space. We extend the so called "backstepping method" by introducing appropriate time-dependent integral transformations in order to map our initial system to a new one which has desired stability properties. The kernels of the integral transformations involved are solutions to non standard multi-dimensional hyperbolic PDEs, where the time dependence introduces several new difficulties in the treatment of their well-posedness. This work generalizes previous results of the literature, where only time-independent systems were considered.


Keywords: Hyperbolic systems, Boundary stabilization, Non-autonomous systems, Backstepping method.

## 1 Introduction and main result

In the present paper we are interested in the one-sided boundary stabilization in finite time of one-dimensional linear hyperbolic balance laws when the coupling coefficients of the system depend on both time and space variables. To investigate this stabilization property we use the by now socalled "backstepping method", a method that consists in transforming our initial system into another system - called target system - for which the stabilization properties are simpler to study. In finite dimension it relies on a recursive design procedure, which in the case of partial differential equations leads to Volterra transformations of the second kind.

The idea of the possibility to transform a control system into another one in order to study its controllability or stabilization properties already goes back to the development of the control theory for linear finite-dimensional systems in the late 60's, notably with the celebrated work Bru70 where the author introduced the so-called "control canonical form". Concerning infinite-dimensional systems, such as systems modeled by partial differential equations (PDEs), this approach is much more complicated. The first attempt in this direction seems to be Rus78, where the author was interested in the spectral determination (i.e. pole placement) of a particular $2 \times 2$ first-order hyperbolic

[^0]system. The difficult task in this approach is, in general, to find an invertible transformation that allows to pass from one system to another and, to the best of our knowledge, there is no general theory for infinite-dimensional systems so far (if possible). In Rus78, the author proposed to use a Volterra transformation of the second kind to pass from what he called the "control normal form" to the control canonical form of his hyperbolic system and, in this way, easily solved his spectral determination problem. In that paper, the use of such a transformation was justified by the analogy with finite-dimensional systems when using transformations of the simple form $\mathrm{Id}+K$ with $K$ being a triangular matrix (while for Volterra transformations of the second kind, $K$ is an integral operator whose kernel is supported in a triangular domain). The use of a Volterra transformation of the second kind to transform a PDE into another one was also introduced at almost the same time in Col77. Therein, the author showed that a one-dimensional perturbed heat equation, with a time and space dependent perturbation, can be transformed into the classical heat equation by means of a Volterra transformation of the second kind whose kernel has to satisfy some PDE posed on a non-standard domain which is triangular. The equation that the kernel has to satisfy is now commonly referred to as the "kernel equation" and the method was then referred by the author of Col77 to as the "method of integral operators". The result of [Col77 was notably applied in [Sei84] to deduce the boundary null-controllability in one space dimension of the perturbed heat equation from that of the classical heat equation.

In the 90 's a method with similar spirit appeared under the name of "backstepping method". This method was primarily designed to transform, thanks to a recursive procedure, finite-dimensional control systems, which may be nonlinear, into control systems which can be stabilized by means of simple feedback laws. This method was later on extended to linear PDEs. The first result in this direction is in CdN98 for a beam equation; see also LK00 for a Burgers' equation. However, the main breakthrough for the PDEs case are in BKL01 BK02, Liu03, which deal with 1-D heat equations and where Volterra transformations of the second kind are introduced or used. In particular in BK02 the backstepping recursive procedure in finite dimension is applied to the semi-discretized finite difference approximation of these equations and it is proved that, as the spatial step size tends to 0 , the backstepping transformation at the finite dimensional level is converging to a Volterra transformation of the second kind. The fact that the transformation which appears with this approach is a Volterra transformation of the second kind comes from the recursive procedure of the backstepping method. With this method the authors, directly inspired by the backstepping in finite dimension, independently arrived at the use of exactly the same transformation as in the two above mentioned pioneering references Rus78 and Col77. This is the reason of the use of the terminology "backstepping" for the construction of stabilizing feedback laws relying on the use of Volterra transformations of the second kind to transform a given control PDE to another control PDE (called the target system) which can be easily stabilized (usually with the null feedback law).

The use of Volterra transformations of the second kind also matches very well with the boundary stabilization of one-dimensional systems since this transformation somehow removes the undesirable terms (or adds desirable ones) of the equation by "bringing" them to the part of the boundary where the feedback is acting (through the kernel equations). This approach rapidly turned out to be very successful in the study of the boundary stabilization of various important PDEs such as heat equations, wave equations, Schrödinger equations, Korteweg-de Vries equations, Kuramoto-Sivashinsky equations, etc. and it eventually leads to the by now reference book [KS08] on this subject. This method is nowadays systematically used as a standard tool to analyze the boundary stabilization for (mainly one-dimensional) PDEs. This method has also received some recent developments. Notably, the use of Volterra transformations of the second kind has started to show some serious limitation for some problems and it has been replaced by more general integral transformations such as Fredholm integral transformations (see e.g. [CL14, CL15, BAK15, CHO16, CHO17, CGM18]) or other kind of integral transformations (see e.g. [SGK09]). In these cases the transformation on the state does not have any special structure and the method is no longer related to the finite dimensional backstepping approach. It is related to the older notion of feedback equivalence, as initiated in Bru70; see also

Kal72, Won85, Section 5.7], and Son98, Section 5.2].
Concerning more specifically systems of hyperbolic equations and the finite-time stabilization property, which is the focus of this article, the first result was obtained in CVKB13. In this paper, the authors developed the original backstepping method to prove the boundary stabilization of a $2 \times 2$ hyperbolic system in finite time, with the best time that can be achieved. The generalization of the result of CVKB13] to $n \times n$ systems was a non-trivial task which was eventually solved in [HDMVK16, HVDMK19] using the ideas introduced previously in HDM15] for $3 \times 3$ systems. The key point was to add additional constraints on the kernel to obtain a specific structure of the coupling parameter in the target system. The time of stabilization found in [HDMVK16, HVDMK19] was then improved in ADM16, CHO17, using two different target systems.

The goal of the present article is to extend the results of the previously mentioned references to time-dependent systems. For the finite-time stabilization of non-autonomous hyperbolic systems, the only works that we are aware of are [DJK16] and AA18 which concerned a single equation with constant speed. Therefore, the non-autonomous case for systems was still left without investigation. The introduction of the time variable in the coupling coefficients obviously complicates the whole situation. As in Col77, DJK16, AA18 we need to introduce integral transformations with timedependent kernels, resulting in much more complex kernel equations to solve. Finally, in addition to the previous references, we would also like to mention the work Wan06 on time-dependent quasilinear hyperbolic systems concerning the related notion of controllability and the works [SK05, KD19], with the references therein, concerning the stabilization of time-dependent parabolic systems (where strong regularity conditions are required to make the backstepping method work, because of the result of Kan90).

The rest of this paper is organized as follows. In the remaining part of Section 11 we present in details the class of hyperbolic systems that we consider and we state our main result. In Section 2 we perform several transformations to show that our initial system can be mapped to a target system which is finite-time stable with desired settling time. In Section 3 we prove the existence and regularity of the kernels of the integral transformations that were used in the previous section. Finally, we gathered in Appendices A, B and C some auxiliary results.

### 1.1 System description

In this article, we focus on the following general $n \times n$ linear hyperbolic systems, which appear for instance in the linearized Saint-Venant equations, plug flow chemical reactors equations, heat exchangers equations and many other physical models of balance laws (see e.g. [BC16, Chapter 1]) around time-varying trajectories:

$$
\left\{\begin{array}{l}
\frac{\partial y}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial y}{\partial x}(t, x)=M(t, x) y(t, x)  \tag{1}\\
y-(t, 1)=u(t), \quad y_{+}(t, 0)=Q(t) y_{-}(t, 0) \\
y\left(t^{0}, x\right)=y^{0}(x)
\end{array}\right.
$$

In (1), $t>t^{0} \geq 0$ and $x \in(0,1), y(t, \cdot)$ is the state at time $t, y^{0}$ is the initial data at time $t^{0}$ and $u(t)$ is the control at time $t$. The matrix $M$ couples the equations of the system inside the domain and the matrix $Q$ couples the equations of the system on the boundary $x=0$. We assume that the matrix $\Lambda$ is diagonal:

$$
\begin{equation*}
\Lambda=\operatorname{diag}\left(\lambda_{1}, \ldots,, \lambda_{n}\right) \tag{2}
\end{equation*}
$$

We denote by $m \in\{1, \ldots, n-1\}$ the number of equations with negative speeds and by $p=n-m \in$ $\{1, \ldots, n-1\}$ the number of equations with positive speeds (all along this work we assume that $n \geq 2$, see Remark 1.11 below for the case $m=n \geq 1$ ). We assume that there exists some $\varepsilon>0$ such
that, for every $t \geq 0$ and $x \in[0,1]$, we have

$$
\begin{equation*}
\lambda_{1}(t, x)<\cdots<\lambda_{m}(t, x)<-\varepsilon<0<\varepsilon<\lambda_{m+1}(t, x)<\cdots<\lambda_{n}(t, x) \tag{3}
\end{equation*}
$$

and, for every $i \in\{1, \ldots, n-1\}$,

$$
\begin{equation*}
\lambda_{i+1}(t, x)-\lambda_{i}(t, x)>\varepsilon \tag{4}
\end{equation*}
$$

Assumptions (3) and (4) will be commented, respectively, in Remarks 1.9 and 1.10 below.
All along this paper, for a vector (or vector-valued function) $v \in \mathbb{R}^{n}$ and a matrix (or matrix-valued function) $A \in \mathbb{R}^{n \times n}$, we use the notation

$$
v=\binom{v_{-}}{v_{+}}, \quad A=\left(\begin{array}{cc}
A_{--} & A_{-+} \\
A_{+-} & A_{++}
\end{array}\right)
$$

where $v_{-} \in \mathbb{R}^{m}, v_{+} \in \mathbb{R}^{p}$ and $A_{--} \in \mathbb{R}^{m \times m}, A_{-+} \in \mathbb{R}^{m \times p}, A_{+-} \in \mathbb{R}^{p \times m}, A_{++} \in \mathbb{R}^{p \times p}$.
We will always assume the following regularities for the parameters involved in the system (1):

$$
\Lambda \in C^{1}([0,+\infty) \times[0,1])^{n \times n}, \quad M \in C^{0}([0,+\infty) \times[0,1])^{n \times n}, \quad Q \in C^{0}([0,+\infty))^{p \times m}
$$

$$
\begin{equation*}
\Lambda, \frac{\partial \Lambda}{\partial x}, M \in L^{\infty}((0,+\infty) \times(0,1))^{n \times n}, \quad Q \in L^{\infty}(0,+\infty)^{p \times m} \tag{5}
\end{equation*}
$$

In this article, we use the notion of "solution along the characteristics" or "broad solution" for the system (1). The necessary background on this notion is given in Appendix A (see also Bre00, Section 3.4] for more information). For the moment we only need to know that, for every $F \in$ $L^{\infty}((0,+\infty) \times(0,1))^{m \times n}, t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$, there exists a unique (broad) solution $y \in$ $C^{0}\left(\left[t^{0},+\infty\right) ; L^{2}(0,1)^{n}\right)$ to the system (1) with

$$
\begin{equation*}
u(t)=\int_{0}^{1} F(t, \xi) y(t, \xi) d \xi \tag{6}
\end{equation*}
$$

The relation (6) will be called the feedback law and the function $F$ will be called the state-feedback gain function.

Let us now give the notion of stability that we are interested in this article (see, for example, [BB98, Definition], BR05, Section 3.2] and Cor07, Definitions 11.11 and 11.27] for time-varying systems in finite dimension).

Definition 1.1. Let $T>0$. We say that the system (1) with feedback law (6) is finite-time stable with settling time $T$ if the following two properties hold:
(i) Finite-time global attractor. For every $t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$,

$$
\begin{equation*}
y\left(t^{0}+T, \cdot\right)=0 \tag{7}
\end{equation*}
$$

(ii) Uniform stability. For every $\varepsilon>0$, there exists $\delta>0$ such that, for every $t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$,

$$
\begin{equation*}
\left(\left\|y^{0}\right\|_{L^{2}(0,1)^{n}} \leq \delta\right) \quad \Longrightarrow \quad\left(\|y(t, \cdot)\|_{L^{2}(0,1)^{n}} \leq \varepsilon, \quad \forall t \geq t^{0}\right) \tag{8}
\end{equation*}
$$

Remark 1.2. The property (8) guarantees that, inside any time interval of the form $\left[t^{0}, t^{0}+T\right]$, the solution is controlled solely by its value at the initial time $t^{0}$, even if this time $t^{0}$ is very large. For
our system (1) this property is in fact a consequence of the first property (7) and that the statefeedback gain function $F$ is in $L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ (see Remark A.3). Such an implication is in general not true for time-dependent hyperbolic systems. A simple example is the following transport equation:

$$
\left\{\begin{array}{l}
\frac{\partial y}{\partial t}(t, x)-\frac{\partial y}{\partial x}(t, x)=0 \\
y(t, 1)=f(t) \int_{0}^{1} y(t, \xi) d \xi \\
y\left(t^{0}, x\right)=y^{0}(x)
\end{array}\right.
$$

where $f \in C^{\infty}([0,+\infty))$ is such that, for every $k \in \mathbb{N}$,

$$
\left\{\begin{array}{l}
f(t)=0, \quad \forall t \in[2 k, 2 k+1] \\
f(t)=t, \quad \forall t \in\left[2 k+\frac{5}{4}, 2 k+\frac{7}{4}\right]
\end{array}\right.
$$

(note that $f \notin L^{\infty}(0,+\infty)$ ). Then the finite-time global attractor property holds (with $T=3$ ) but the uniform stability property does not hold (consider the sequences $y_{\delta}^{0}(x)=\delta$ for every $x \in(0,1)$ and $t_{\delta}^{0}=2\left\lceil\frac{1}{\delta}\right\rceil+\frac{5}{4}$, where $\lceil\cdot\rceil$ denotes the ceiling function).
Remark 1.3. As we are trying to find a state-feedback gain function $F$ so that (1) with feedback law (6) is finite-time stable, let us first point out that, in general, $F=0$ does not work. A simple example is provided by the $2 \times 2$ system with constant coefficients ( $t^{0}=0$ to simplify)

$$
\left\{\begin{array}{l}
\frac{\partial y_{-}}{\partial t}(t, x)-\frac{\partial y_{-}}{\partial x}(t, x)=-c y_{+}(t, x) \\
\frac{\partial y_{+}}{\partial t}(t, x)+\frac{\partial y_{+}}{\partial x}(t, x)=-c y_{-}(t, x) \\
y_{-}(t, 1)=0, \quad y_{+}(t, 0)=y_{-}(t, 0) \\
y(0, x)=y^{0}(x)
\end{array}\right.
$$

which is exponentially unstable for $c>\pi$ (see e.g. BC16, Proposition 5.12] with $y_{-}(t, x)=S_{1}(t, 1-x)$ and $y_{+}(t, x)=S_{2}(t, 1-x)$ ), and thus not finite-time stable.

### 1.2 The characteristics

To state the main result of this paper we need to introduce the characteristic curves associated with system (1). To this end, it is convenient to first extend $\Lambda$ to a function of $\mathbb{R}^{2}$ (still denoted by $\Lambda$ ).
Remark 1.4. This extension procedure can be done in such a way that the properties (2), (3), (4) and (5) remain valid on $\mathbb{R}^{2}$. We can take for instance

$$
\bar{\lambda}_{i}(t, x)= \begin{cases}\lambda_{i}(t, x) & \text { if } t \geq 0 \\ \lambda_{i}(0, x)+\delta\left(\lambda_{i}(0, x)-\lambda_{i}\left(1-e^{t / \delta}, x\right)\right) & \text { if } t<0\end{cases}
$$

where $\delta>0$ is small enough so that $-\varepsilon+4 \delta \max _{i}\left\|\lambda_{i}\right\|_{L^{\infty}((0,1) \times(0,1))}<-\varepsilon / 2$ to guarantee the properties (3) and (4) with $\varepsilon / 2$ in place of $\varepsilon$. This extends the function to $\mathbb{R} \times[0,1]$. We can use a similar procedure to then extend it to $\mathbb{R}^{2}$. We can check that the results of this paper do not depend on such a choice of extension (all the important data are uniquely determined on the domain of interest $(0,+\infty) \times(0,1))$.

### 1.2.1 The flow

For every $i \in\{1, \cdots, n\}$, let $\chi_{i}$ be the flow associated with $\lambda_{i}$, i.e. for every $(t, x) \in \mathbb{R} \times \mathbb{R}$, the function $s \mapsto \chi_{i}(s ; t, x)$ is the solution to the ODE

$$
\left\{\begin{array}{l}
\frac{\partial \chi_{i}}{\partial s}(s ; t, x)=\lambda_{i}\left(s, \chi_{i}(s ; t, x)\right), \quad \forall s \in \mathbb{R}  \tag{9}\\
\chi_{i}(t ; t, x)=x
\end{array}\right.
$$

The existence and uniqueness of the solution to the ODE (9) follows from the (local) Cauchy-Lipschitz theorem and this solution is global since $\lambda_{i}$ is bounded (by the finite time blow-up theorem, see e.g. Har02, Theorem II.3.1]). The uniqueness of the solution to the ODE (9) also yields the group property

$$
\begin{equation*}
\chi_{i}\left(\sigma ; s, \chi_{i}(s ; t, x)\right)=\chi_{i}(\sigma ; t, x), \quad \forall \sigma \in \mathbb{R} \tag{10}
\end{equation*}
$$

By classical regularity results on ODEs (see e.g. Har02, Theorem V.3.1]), $\chi_{i}$ has the regularity

$$
\begin{equation*}
\chi_{i} \in C^{1}\left(\mathbb{R}^{3}\right) \tag{11}
\end{equation*}
$$

and, for every $s, t, x \in \mathbb{R}$, we have

$$
\begin{equation*}
\frac{\partial \chi_{i}}{\partial t}(s ; t, x)=-\lambda_{i}(t, x) e^{\int_{t}^{s} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta}, \quad \frac{\partial \chi_{i}}{\partial x}(s ; t, x)=e^{\int_{t}^{s} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta} \tag{12}
\end{equation*}
$$

Note in particular that

$$
\begin{cases}\frac{\partial \chi_{i}}{\partial t}(s ; t, x)>0 & \text { if } i \in\{1, \ldots, m\}  \tag{13}\\ \frac{\partial \chi_{i}}{\partial t}(s ; t, x)<0 & \text { if } i \in\{m+1, \ldots, n\} \\ \frac{\partial \chi_{i}}{\partial x}(s ; t, x)>0 & \end{cases}
$$

### 1.2.2 The entry and exit times

For every $i \in\{1, \cdots, n\}, t \in \mathbb{R}$ and $x \in[0,1]$, let $s_{i}^{\text {in }}(t, x), s_{i}^{\text {out }}(t, x) \in \mathbb{R}$ be the entry and exit times of the flow $\chi_{i}(\cdot ; t, x)$ inside the domain $[0,1]$, i.e. the respective unique solutions to

$$
\left\{\begin{array}{ll}
\chi_{i}\left(s_{i}^{\text {in }}(t, x) ; t, x\right)=1, & \chi_{i}\left(s_{i}^{\text {out }}(t, x) ; t, x\right)=0,  \tag{14}\\
\chi_{i}\left(s_{i}^{\text {in }}(t, x) ; t, x\right)=0, & \chi_{i}\left(s_{i}^{\text {out }}(t, x) ; t, x\right)=1,
\end{array} \quad \text { if } i \in\{1, \ldots, m\}, 1, \ldots, n\right\} .
$$

The existence and uniqueness of $s_{i}^{\text {out }}(t, x)$ and $s_{i}^{\text {in }}(t, x)$ are guaranteed by the assumption (3). Note that we always have

$$
\begin{equation*}
s_{i}^{\text {in }}(t, x) \leq t \leq s_{i}^{\text {out }}(t, x) \tag{15}
\end{equation*}
$$

and the cases of equalities are given by

$$
\left\{\begin{array}{rlll}
\mathbf{1 8 7}^{\text {in }}(t, x)=t & \Longleftrightarrow x=1, & s_{i}^{\text {out }}(t, x)=t & \Longleftrightarrow  \tag{16}\\
s_{i}^{\text {in }} & \Longleftrightarrow x=0, & \text { if } i \in\{1, \ldots, m\} \\
s_{i}^{\text {in }}(t, x)=t & \Longleftrightarrow x=0, & s_{i}^{\text {out }}(t, x)=t & \Longleftrightarrow
\end{array}\right.
$$

It readily follows from (10) and the uniqueness of $s_{i}^{\text {in }}, s_{i}^{\text {out }}$ that, for every $s \in\left[s_{i}^{\text {in }}(t, x), s_{i}^{\text {out }}(t, x)\right]$,

$$
\begin{equation*}
s_{i}^{\text {in }}\left(s, \chi_{i}(s ; t, x)\right)=s_{i}^{\text {in }}(t, x), \quad s_{i}^{\text {out }}\left(s, \chi_{i}(s ; t, x)\right)=s_{i}^{\text {out }}(t, x) . \tag{17}
\end{equation*}
$$

From (11) and by the implicit function theorem, we have

$$
\begin{equation*}
s_{i}^{\text {in }}, s_{i}^{\text {out }} \in C^{1}(\mathbb{R} \times[0,1]) \tag{18}
\end{equation*}
$$

Moreover, integrating the ODE (9) and using the assumption (3), we have the following bounds, valid for every $t \in \mathbb{R}$ and $x \in[0,1]$,

$$
\begin{equation*}
t-s_{i}^{\text {in }}(t, x)<\frac{1}{\varepsilon}, \quad s_{i}^{\text {out }}(t, x)-t<\frac{1}{\varepsilon} . \tag{19}
\end{equation*}
$$

On the other hand, differentiating (14) and using (13) with (3), we see that, for every $t \in \mathbb{R}$ and $x \in[0,1]$, we have

$$
\left\{\begin{array}{cll}
\frac{\partial s_{i}^{\text {in }}}{\partial t}(t, x)>0, & \frac{\partial s_{i}^{\text {out }}}{\partial t}(t, x)>0, &  \tag{20}\\
\frac{\partial s_{i}^{\text {in }}}{\partial x}(t, x)>0, & \frac{\partial s_{i}^{\text {out }}}{\partial x}(t, x)>0 & \text { if } i \in\{1, \ldots, m\}, \\
\frac{\partial s_{i}^{\text {in }}}{\partial x}(t, x)<0, & \frac{\partial s_{i}^{\text {out }}}{\partial x}(t, x)<0 & \text { if } i \in\{m+1, \ldots, n\} .
\end{array}\right.
$$

Finally, from the assumption (3) and classical results on comparison for ODEs (see e.g. Har02, Corollary III.4.2]), we have, for every $t \in \mathbb{R}$ and $x \in[0,1]$,

$$
\left\{\begin{align*}
s_{m}^{\text {in }}(t, x)<\ldots<s_{1}^{\text {in }}(t, x) & \text { if } x \neq 1, \quad s_{1}^{\text {out }}(t, x)<\ldots<s_{m}^{\text {out }}(t, x) \quad \text { if } x \neq 0  \tag{21}\\
s_{m+1}^{\text {in }}(t, x)<\ldots<s_{n}^{\text {in }}(t, x) & \text { if } x \neq 0, \quad s_{n}^{\text {out }}(t, x)<\ldots<s_{m+1}^{\text {out }}(t, x) \quad \text { if } x \neq 1
\end{align*}\right.
$$

### 1.3 Main result and comments

We are now in position to state the main result of this paper:
Theorem 1.5. Let $\Lambda, M$ and $Q$ satisfy (2), (3), (4) and (5). Then, there exists a state-feedback gain function $F \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ such that the system (1) with feedback law (6) is finite-time stable with settling time $T_{\text {unif }}(\Lambda)$ defined by

$$
\begin{equation*}
T_{\text {unif }}(\Lambda)=\sup _{t^{0} \geq 0} s_{m+1}^{\text {out }}\left(s_{m}^{\text {out }}\left(t^{0}, 1\right), 0\right)-t^{0} \tag{22}
\end{equation*}
$$

Moreover, if for some $\tau>0, \Lambda, M$ and $Q$ are $\tau$-periodic with respect to time (that is $\Lambda(t+\tau, x)=$ $\Lambda(t, x)$ for every $t \geq 0$ and $x \in[0,1]$, same for $M$ and $Q$ ) then one can also impose to $F$ to be $\tau$-periodic with respect to time (almost everywhere).

Let us remark that, thanks to (15), (16) and (19), we always have

$$
0<T_{\text {unif }}(\Lambda)<\frac{2}{\varepsilon}
$$

Note as well, thanks to (21) and the first line in (20), that we have

$$
T_{\text {unif }}(\Lambda)=\max _{j \in\{m+1, \ldots, n\}} \max _{i \in\{1, \ldots, m\}} \sup _{t^{0} \geq 0} s_{j}^{\text {out }}\left(s_{i}^{\text {out }}\left(t^{0}, 1\right), 0\right)-t^{0} .
$$

Example 1.6. Theorem 1.5 applies for instance to the following coupled $2 \times 2$ system:

$$
\left\{\begin{array}{l}
\frac{\partial y_{1}}{\partial t}(t, x)-\frac{\partial y_{1}}{\partial x}(t, x)=m_{11}(t, x) y_{1}(t, x)+m_{12}(t, x) y_{2}(t, x),  \tag{23}\\
\frac{\partial y_{2}}{\partial t}(t, x)+\left(1+\frac{1}{1+t}\right) \frac{\partial y_{2}}{\partial x}(t, x)=m_{21}(t, x) y_{1}(t, x)+m_{22}(t, x) y_{2}(t, x), \\
y_{1}(t, 1)=u(t), \quad y_{2}(t, 0)=q(t) y_{1}(t, 0), \\
y\left(t^{0}, x\right)=y^{0}(x),
\end{array}\right.
$$

where $M=\left(m_{i j}\right)_{1 \leq i, j \leq 2}$ and $Q=(q)$ are any parameters with the regularity (5). Let us show how to compute $T_{\text {unif }}(\bar{\Lambda})$ for this example. First of all, it is clear that $\chi_{1}(s ; t, x)=-s+t+x$, so that $s_{1}^{\text {out }}\left(t^{0}, 1\right)=t^{0}+1$. On the other hand, we have $\chi_{2}(s ; t, x)=s+\ln (1+s)-t-\ln (1+t)+x$. Therefore, $h\left(t^{0}\right)=s_{2}^{\text {out }}\left(t^{0}+1,0\right)-t^{0}$ solves $\Psi\left(h\left(t^{0}\right), t^{0}\right)=0$, where

$$
\Psi\left(h, t^{0}\right)=h+\ln \left(1+h+t^{0}\right)-2-\ln \left(2+t^{0}\right) .
$$

Taking the derivative of the relation $\Psi\left(h\left(t^{0}\right), t^{0}\right)=0$ and using the fact that $h \geq 1$ by (15), we see that $h^{\prime}\left(t^{0}\right) \geq 0$, so that $h$ is non-decreasing. Since $h \leq 2$ by (19), the function $h$ is thus a bounded nondecreasing function and, consequently, $\lim _{t^{0} \rightarrow+\infty} h\left(t^{0}\right)$ exists and is equal to $\sup _{t^{0} \geq 0} h\left(t^{0}\right)=T_{\text {unif }}(\Lambda)$. Writing the relation $\Psi\left(h\left(t^{0}\right), t^{0}\right)=0$ as follows for $t^{0}>0$

$$
h\left(t^{0}\right)+\ln \left(\frac{1}{t^{0}}+\frac{h\left(t^{0}\right)}{t^{0}}+1\right)-2-\ln \left(\frac{2}{t^{0}}+1\right)=0,
$$

and letting $t^{0} \rightarrow+\infty$ we obtain the value $T_{\text {unif }}(\Lambda)=2$.
Remark 1.7. Observe that the time $T_{\text {unif }}(\Lambda)$ does not depend on the parameters $M$ and $Q$. It depends only on $\Lambda$ on $[0,+\infty) \times(0,1)$. Moreover, this is the best time one can obtain, uniformly with respect to all the possible choices of $M$ and $Q$ (this explains our notation " $T_{\text {unif }}(\Lambda)$ "). More precisely,

$$
T_{\text {unif }}(\Lambda)=\min E,
$$

where $E$ is the set of $T>0$ such that, for every $M$ and $Q$ with the regularity (5), there exists a state-feedback gain function $F \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ so that the system (1] with feedback law (6) is finite-time stable with settling time $T$. Indeed, Theorem 1.5 establishes that $T_{\text {unif }}(\Lambda) \in E$, so that $E \neq \emptyset$. On the other hand, taking $M=0$ and the constant matrix

we can check from the very definition of broad solution (see Definition A.1) that, if $T<T_{\text {unif }}(\Lambda)$, then there exist $t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$ such that the corresponding solution to (1) satisfies $y\left(t^{0}+T, \cdot\right) \neq 0$, whatever $u \in L^{\infty}\left(t^{0}, t^{0}+T\right)^{m}$ is.

Of course, for particular choices of $M$ and $Q$ one may obtain a better settling time (a trivial example being $M=0$ and $Q=0$ ). In the case of time-independent systems, the minimal time in which one can achieve the stabilization and related controllability properties has been recently discussed in [CN19] and HO19 (see also the references therein).

Remark 1.8. If the speeds do not depend on time, i.e. $\lambda_{\ell}(t, x)=\lambda_{\ell}(x)$ for every $t \geq 0$, then we have a more explicit formula for the time $T_{\text {unif }}(\Lambda)$, namely:

$$
\begin{equation*}
T_{\text {unif }}(\Lambda)=\int_{0}^{1} \frac{1}{-\lambda_{m}(\xi)} d \xi+\int_{0}^{1} \frac{1}{\lambda_{m+1}(\xi)} d \xi \tag{24}
\end{equation*}
$$

The value 24 is obtained by integrating over $\xi \in[0,1]$ the differential equation satisfied by the inverse functions $\xi \longmapsto \chi_{m}^{-1}(\xi ; t, 1)$ and $\xi \longmapsto \chi_{m+1}^{-1}(\xi ; t, 0)$.
Remark 1.9. The assumption (3) that the negative (resp. positive) speeds are uniformly bounded from above (resp. below), despite not being necessary for the existence of a solution to (1), is to be expected for the system (1) to be finite-time stable. This is an issue that is not specific to systems and that already occurs for a single equation. Indeed, let us consider for instance the equation with speed $\lambda(t)=-e^{-t}$ (and $t^{0}=0$ to simplify):

$$
\left\{\begin{array}{l}
\frac{\partial y}{\partial t}(t, x)-e^{-t} \frac{\partial y}{\partial x}(t, x)=0 \\
y(t, 1)=u(t) \\
y(0, x)=y^{0}(x)
\end{array}\right.
$$

Then, whatever $y^{0} \in L^{2}(0,1)$ and $u \in L^{\infty}(0,+\infty)$ are, if $y^{0} \neq 0$ in a neighborhood of 1 we have

$$
y(T, \cdot) \neq 0, \quad \forall T>0
$$

This is easily seen thanks to the explicit representation of the solution (obtained by the characteristic method):

$$
y(t, x)=\left\{\begin{array}{cl}
y^{0}\left(1-\left(e^{-t}-x\right)\right) & \text { if } 0<x<e^{-t} \\
u\left(\ln \left(\frac{1}{1+e^{-t}-x}\right)\right) & \text { if } e^{-t}<x<1
\end{array}\right.
$$

Remark 1.10. Contrary to (3), the assumption (4) is mainly technical. This assumption is needed because we will have to divide in the sequel by the quantities $\lambda_{j}-\lambda_{i}$ (see in particular (67) below) and we will need this inverse function to be bounded. However, this condition is clearly not necessary for some systems (1) to be finite-time stable. Indeed, consider for instance the following $3 \times 3$ system:

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$$
\left\{\begin{array}{l}
\frac{\partial y_{1}}{\partial t}(t, x)-\frac{\partial y_{1}}{\partial x}(t, x)=y_{2}(t, x)  \tag{25}\\
\frac{\partial y_{2}}{\partial t}(t, x)-\left(1-\frac{e^{-t}}{2}\right) \frac{\partial y_{2}}{\partial x}(t, x)=0 \\
\frac{\partial y_{3}}{\partial t}(t, x)+\frac{\partial y_{3}}{\partial x}(t, x)=0 \\
y_{1}(t, 1)=u_{1}(t), \quad y_{2}(t, 1)=u_{2}(t), \quad y_{3}(t, 0)=y_{2}(t, 0) \\
y\left(t^{0}, x\right)=y^{0}(x)
\end{array}\right.
$$

Then, using the characteristic method it is not difficult to see that the system with $u_{1}=u_{2}=0$ is finite-time stable with settling time $T+1$, where $T$ is the unique positive solution to the equation $T+\frac{e^{-T}}{2}=\frac{3}{2}$.
Remark 1.11. The case $m=n \geq 1$ (no boundary conditions at $x=0$ ) is easier and does not require the techniques presented in this paper. Indeed, it can be checked using for instance the constructive method of LR03, Wan06 that in this case the system (1) with $u=0$ is finite-time stable with settling time equal to $\sup _{t^{0} \geq 0} s_{m}^{\text {out }}\left(t^{0}, 1\right)-t^{0}$.
e Section 3 below for the sake of the presentation.

Finally, because of the nature of the transformations that we will use in the sequel, we are led to consider a class of systems that is slightly more general than (1). All the systems of this paper will have the following form:

$$
\left\{\begin{array}{l}
\frac{\partial y}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial y}{\partial x}(t, x)=M(t, x) y(t, x)+G(t, x) y(t, 0),  \tag{26}\\
y_{-}(t, 1)=\int_{0}^{1} F(t, \xi) y(t, \xi) d \xi, \quad y_{+}(t, 0)=Q(t) y_{-}(t, 0), \\
y\left(t^{0}, x\right)=y^{0}(x),
\end{array}\right.
$$

where $M$ and $Q$ will have at least the regularity (5), $F \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ and

$$
G \in C^{0}([0,+\infty) \times[0,1])^{n \times n} \cap L^{\infty}((0,+\infty) \times(0,1))^{n \times n}
$$

Therefore, (26) is similar to (1) but has the extra term with $G$. In what follows, we will also refer to a system of the form (26) as

$$
(M, G, F, Q)
$$

Hyperbolic equations similar to $(0, G, F, Q)$ were called in "control normal form" in the pioneering work [Rus78, p. 212] for the similarity with the finite-dimensional setting (see also the earlier paper (Bru70).

### 2.1 Removal of the diagonal terms

In this section we just perform a simple preliminary transformation in order to remove the diagonal terms in $M$. This is only a technical step, which is nevertheless necessary in view of the existence of the transformation that we will use in the next section, see Remark 2.7 below. This step is sometimes called "exponential pre-transformation" in the case of time-independent systems (see Remark 2.3 below). More precisely, the goal of this section is to establish the following result:

Proposition 2.1. There exists $M^{1}=\left(m_{i j}^{1}\right)_{1 \leq i, j \leq n} \in C^{0}([0,+\infty) \times[0,1])^{n \times n} \cap L^{\infty}((0,+\infty) \times$ $(0,1))^{n \times n}$ with diagonal terms equal to zero:

$$
\begin{equation*}
m_{i i}^{1}=0, \quad \forall i \in\{1, \ldots, n\}, \tag{27}
\end{equation*}
$$

and there exists $Q^{1} \in C^{0}([0,+\infty))^{p \times m} \cap L^{\infty}(0,+\infty)^{p \times m}$ such that, for every $F^{1} \in L^{\infty}((0,+\infty) \times$ $(0,1))^{m \times n}$, there exists $F \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ such that the following property holds for every $T>0$ :
$\left(M^{1}, 0, F^{1}, Q^{1}\right)$ is finite-time stable with settling time $T$
$\Longrightarrow \quad(M, 0, F, Q)$ is finite-time stable with settling time $T$.

### 2.1.1 Formal computations

To prove Proposition 2.1 the idea is to show that, for every $F^{1}$, there exists $F$ such that we can transform a solution of $(M, 0, F, Q)$ into a solution of $\left(M^{1}, 0, F^{1}, Q^{1}\right)$. Let then $y$ be the solution to the system $(M, 0, F, Q)$ with state-feedback gain function $F$ to be determined below and initial data $y^{0}$. Let $\Phi:[0,+\infty) \times[0,1] \longrightarrow \mathbb{R}^{n \times n}$ be a smooth matrix-valued function and set

$$
\begin{equation*}
w(t, x)=\Phi(t, x) y(t, x) \tag{29}
\end{equation*}
$$

Let us now perform some formal computations in order to see what $w$ can solve. Using the equation satisfied by $y$, we have

$$
\frac{\partial w}{\partial t}+\Lambda \frac{\partial w}{\partial x}=\left(\frac{\partial \Phi}{\partial t}+\Phi M+\Lambda \frac{\partial \Phi}{\partial x}\right) y+(-\Phi \Lambda+\Lambda \Phi) \frac{\partial y}{\partial x}
$$

On the other hand, using the boundary condition satisfied by $y$ at $x=0$, we have

$$
\begin{aligned}
w_{+}(t, 0)-Q^{1}(t) & w_{-}(t, 0) \\
& =\left(\Phi_{+-}(t, 0)+\Phi_{++}(t, 0) Q(t)-Q^{1}(t) \Phi_{--}(t, 0)-Q^{1}(t) \Phi_{-+}(t, 0) Q(t)\right) y_{-}(t, 0)
\end{aligned}
$$

Finally, at $x=1$, we have
$w_{-}(t, 1)-\int_{0}^{1} F^{1}(t, \xi) w(t, \xi) d \xi=\int_{0}^{1}\left(\Phi_{--}(t, 1) F(t, \xi)-F^{1}(t, \xi) \Phi(t, \xi)\right) y(t, \xi) d \xi+\Phi_{-+}(t, 1) y_{+}(t, 1)$.
Thus, we see that $w$ satisfies at $x=1$ the boundary condition $w_{-}(t, 1)=\int_{0}^{1} F^{1}(t, \xi) w(t, \xi) d \xi$ if $\Phi_{-+}(t, 1)=0$ and

$$
\begin{equation*}
F(t, \xi)=\Phi_{--}(t, 1)^{-1} F^{1}(t, \xi) \Phi(t, \xi) \tag{30}
\end{equation*}
$$

provided that $\Phi_{--}(t, 1)$ is also invertible. Moreover, note that $F$ belongs to $L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ provided that $F^{1}$ belongs to this space as well and

$$
\begin{equation*}
\exists C>0, \quad\left\|\Phi_{--}(\cdot, 1)^{-1}\right\|_{L^{\infty}(0,+\infty)^{m \times m}} \leq C \tag{31}
\end{equation*}
$$

In summary, $w$ defined by 29 is the solution of $\left(M^{1}, 0, F^{1}, Q^{1}\right)$ with state-feedback gain function $F^{1}$ (which is assumed to be known) and initial data $w^{0}(\cdot)=\Phi(0, \cdot) y^{0}(\cdot)$ if we have the following four properties:
(i) $\Lambda(t, x) \Phi(t, x)=\Phi(t, x) \Lambda(t, x)$ for every $t \geq 0$ and $x \in[0,1]$.
(ii) The matrices $\Phi(t, x)$ and $\Phi_{--}(t, 0)+\Phi_{-+}(t, 0) Q(t)$ are invertible for every $t \geq 0$ and $x \in[0,1]$.
(iii) $\Phi_{-+}(t, 1)=0$ for every $t \geq 0$ (it then follows with (ii) that $\Phi_{--}(t, 1)$ is invertible).
(iv) $M^{1}$ and $Q^{1}$ are defined by

$$
\left\{\begin{array}{l}
M^{1}(t, x)=\left(\frac{\partial \Phi}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \Phi}{\partial x}(t, x)+\Phi(t, x) M(t, x)\right) \Phi(t, x)^{-1}  \tag{32}\\
Q^{1}(t)=\left(\Phi_{+-}(t, 0)+\Phi_{++}(t, 0) Q(t)\right)\left(\Phi_{--}(t, 0)+\Phi_{-+}(t, 0) Q(t)\right)^{-1}
\end{array}\right.
$$

Finally, it is not difficult to check that the stability property $(28)$ is indeed satisfied since the state-feedback gain function $F$ is solely determined by the state-feedback gain function $F^{1}$ and, at every fixed $t \geq 0$, the transformation 29 defines an injective (in fact, invertible) map of $L^{2}(0,1)^{n}$.

### 2.1.2 Existence of the transformation

Let us now prove the existence of a function $\Phi$ with the properties listed above and which in addition ensures that the condition (27) on $M^{1}$ holds.

Proposition 2.2. There exists $\Phi$ with $\Phi, \frac{\partial \Phi}{\partial t}+\Lambda \frac{\partial \Phi}{\partial x} \in C^{0}([0,+\infty) \times[0,1])^{n \times n} \cap L^{\infty}((0,+\infty) \times$ $(0,1))^{n \times n}$ such that the properties (i), (ii), (iii) and (31) are satisfied and such that the matrixvalued function $M^{1}$ defined in (32) satisfies (27).

Proof. Let $\Phi$ be the diagonal matrix-valued function defined for every $t \geq 0$ and $x \in[0,1]$ by

$$
\Phi(t, x)=\operatorname{diag}\left(\phi_{1}(t, x), \ldots, \phi_{n}(t, x)\right)
$$

where, for every $i \in\{1, \ldots, n\}$,

$$
\begin{equation*}
\phi_{i}(t, x)=e^{-\int_{s_{i}^{\mathrm{in}}(t, x)}^{t} m_{i i}\left(\sigma, \chi_{i}(\sigma ; t, x)\right) d \sigma} \tag{33}
\end{equation*}
$$

where $m_{i i}$ is extended to negative times by an arbitrary function that keeps the regularity (5). Clearly, $\phi_{i} \in C^{0}([0,+\infty) \times[0,1])$ and it follows from (19) that $\phi_{i} \in L^{\infty}((0,+\infty) \times(0,1))$.

It is clear that the first property (i) holds since $\Lambda$ and $\Phi$ are both diagonal matrices. Since $\Phi_{-+}=0$, the third property (iii) is automatically satisfied. It also follows that, to check the second property (ii), we only need to show that $\Phi(t, x)$ is invertible, which readily follows from the explicit expression of $\phi_{i}$. The estimate (31) is obviously true since $\Phi_{--}(t, 1)=\operatorname{Id}_{\mathbb{R}^{m \times m}}$ (recall (16) . Finally, $M^{1}$ defined in (32) satisfies (27) since $\phi_{i}$ satisfies the following linear hyperbolic equation:

$$
\frac{\partial \phi_{i}}{\partial t}(t, x)+\lambda_{i}(t, x) \frac{\partial \phi_{i}}{\partial x}(t, x)+m_{i i}(t, x) \phi_{i}(t, x)=0
$$

Remark 2.3. There are obviously other possible choices for $\phi_{i}$, for instance in the time-independent case we can take the slightly simpler function $\phi_{i}(t, x)=e^{-\int_{0}^{x} \frac{m_{i i}(\xi)}{\lambda_{i}(\xi)} d \xi}$ (which coincides with (33) only for $i \in\{m+1, \ldots, n\})$.

### 2.2 Volterra transformation

In this section we perform a second transformation to remove some coupling terms of the system. The system will then have a triangular coupling structure, which is the key point to show later on (Section 2.3 below) that this system is finite-time stable with settling time $T_{\text {unif }}(\Lambda)$. More precisely, the goal of this section is to establish the following result:
where

$$
\left\{\begin{array}{l}
T_{1}\left(t^{0}\right)=s_{1}^{\text {out }}\left(t^{0}, 1\right), \\
T_{i}\left(t^{0}\right)=s_{i}^{\text {out }}\left(T_{i-1}\left(t^{0}\right), 1\right), \quad \forall i \in\{2, \ldots, m\}
\end{array}\right.
$$

We do not detail this point here because it is not needed, and we refer to the arguments used in the proof of Proposition 2.12 below for an idea of the proof of this assertion. As a result, the combination of Proposition 2.4 with Proposition 2.1 already shows that our initial system $(M, 0, F, Q)$ is finite-time stable for some $F$, with settling time $T\left(0, G^{2}, 0, Q^{1}\right)$. However, this time $T\left(0, G^{2}, 0, Q^{1}\right)$ is always strictly larger than the time $T_{\text {unif }}(\Lambda)$ given in Theorem 1.5 (as long as $m>1$ ). In the case of timeindependent systems, the time $T\left(0, G^{2}, 0, Q^{1}\right)$ is the time obtained in HDMVK16 HVDMK19 and it has the more explicit expression

$$
T\left(0, G^{2}, 0, Q^{1}\right)=\sum_{i=1}^{m} \int_{0}^{1} \frac{1}{-\lambda_{i}(\xi)} d \xi+\int_{0}^{1} \frac{1}{\lambda_{m+1}(\xi)} d \xi
$$

### 2.2.1 Formal computations

Let us now show how to establish Proposition 2.4 As before, the goal is to show that, for every $F^{2}$, there exists $F^{1}$ such that we can transform a solution of ( $M^{1}, 0, F^{1}, Q^{1}$ ) into a solution of $\left(0, G^{2}, F^{2}, Q^{1}\right)$. Let then $w$ be the solution to the system $\left(M^{1}, 0, F^{1}, Q^{1}\right)$ with state-feedback gain function $F^{1}$ to be determined below and initial data $w^{0}$. Inspired by the works mentioned at the beginning of Section 2 we use a Volterra transformation of the second kind as follows:

$$
\begin{equation*}
\gamma(t, x)=w(t, x)-\int_{0}^{x} K(t, x, \xi) w(t, \xi) d \xi \tag{37}
\end{equation*}
$$

where we suppose for the moment that the kernel $K$ is smooth on $\overline{\mathcal{T}}$, where $\mathcal{T}$ is the infinite triangular prism defined by

$$
\mathcal{T}=\{(t, x, \xi) \in(0,+\infty) \times(0,1) \times(0,1), \quad \xi<x\}
$$

Let us now perform some formal computations to see what $\gamma$ can solve. We have

$$
\begin{aligned}
& \frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=\frac{\partial w}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial w}{\partial x}(t, x) \\
& -\int_{0}^{x} \frac{\partial K}{\partial t}(t, x, \xi) w(t, \xi) d \xi-\int_{0}^{x} K(t, x, \xi) \frac{\partial w}{\partial t}(t, \xi) d \xi \\
& -\Lambda(t, x) K(t, x, x) w(t, x)-\Lambda(t, x) \int_{0}^{x} \frac{\partial K}{\partial x}(t, x, \xi) w(t, \xi) d \xi .
\end{aligned}
$$

Using the equation satisfied by $w$, we obtain

$$
\begin{aligned}
& \frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=M^{1}(t, x) w(t, x) \\
&- \int_{0}^{x} \frac{\partial K}{\partial t}(t, x, \xi) w(t, \xi) d \xi- \\
& \quad \int_{0}^{x} K(t, x, \xi)\left(-\Lambda(t, \xi) \frac{\partial w}{\partial \xi}(t, \xi)+M^{1}(t, \xi) w(t, \xi)\right) d \xi \\
&-\Lambda(t, x) K(t, x, x) w(t, x)-\Lambda(t, x) \int_{0}^{x} \frac{\partial K}{\partial x}(t, x, \xi) w(t, \xi) d \xi
\end{aligned}
$$

Integrating by parts the third term of the right hand side and using the boundary condition $w_{+}(t, 0)=$ $Q^{1}(t) w_{-}(t, 0)$, we finally obtain

$$
\begin{aligned}
& \frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=\int_{0}^{x}\left(-\frac{\partial K}{\partial t}(t, x, \xi)-\frac{\partial K}{\partial \xi}(t, x, \xi) \Lambda(t, \xi)-K(t, x, \xi) \frac{\partial \Lambda}{\partial \xi}(t, \xi)\right. \\
& \left.-K(t, x, \xi) M^{1}(t, \xi)-\Lambda(t, x) \frac{\partial K}{\partial x}(t, x, \xi)\right) w(t, \xi) d \xi \\
& +\left(M^{1}(t, x)+K(t, x, x) \Lambda(t, x)-\Lambda(t, x) K(t, x, x)\right) w(t, x)-K(t, x, 0) \Lambda(t, 0)\binom{\operatorname{Id}_{\mathbb{R}^{m \times m}}}{Q^{1}(t)} w_{-}(t, 0)
\end{aligned}
$$

On the other hand, since $\gamma(t, 0)=w(t, 0), \gamma$ satisfies the same boundary condition as $w$ at $x=0$ :

$$
\gamma_{+}(t, 0)-Q^{1}(t) \gamma_{-}(t, 0)=w_{+}(t, 0)-Q^{1}(t) w_{-}(t, 0)=0
$$

Finally, at $x=1$, we have

$$
\begin{aligned}
& \gamma_{-}(t, 1)-\int_{0}^{1} F^{2}(t, \xi) \gamma(t, \xi) d \xi= \\
& \quad \int_{0}^{1}\left(F^{1}(t, \xi)-K_{-}(t, 1, \xi)-F^{2}(t, \xi)+\int_{\xi}^{1} F^{2}(t, \zeta) K(t, \zeta, \xi) d \zeta\right) w(t, \xi) d \xi
\end{aligned}
$$ where $K_{-}$denotes the $m \times n$ sub-matrix of $K$ formed by its first $m$ rows. Thus, we see that $\gamma$ satisfies at $x=1$ the boundary condition $\gamma_{-}(t, 1)=\int_{0}^{1} F^{2}(t, \xi) \gamma(t, \xi) d \xi$ if we take

$$
\begin{equation*}
F^{1}(t, \xi)=K_{-}(t, 1, \xi)+F^{2}(t, \xi)-\int_{\xi}^{1} F^{2}(t, \zeta) K(t, \zeta, \xi) d \zeta \tag{38}
\end{equation*}
$$

Note that $F^{1}$ belongs to $L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ provided that $F^{2}$ belongs to this space as well and

$$
K \in L^{\infty}(\mathcal{T})^{n \times n}
$$

In summary, $\gamma$ defined by $(37)$ is the solution to $\left(0, G^{2}, F^{2}, Q^{1}\right)$ with initial data $\gamma^{0}(x)=w^{0}(x)-$ $\int_{0}^{x} K\left(t^{0}, x, \xi\right) w^{0}(\xi) d \xi$ if we have the following two properties:
(i) For every $(t, x, \xi) \in \mathcal{T}$,

$$
\left\{\begin{array}{l}
\begin{array}{l}
\frac{\partial K}{\partial t}(t, x, \xi)+\Lambda(t, x) \frac{\partial K}{\partial x}(t, x, \xi)+\frac{\partial K}{\partial \xi}(t, x, \xi) \Lambda(t, \xi) \\
\\
\quad+K(t, x, \xi)\left(\frac{\partial \Lambda}{\partial \xi}(t, \xi)+M^{1}(t, \xi)\right)=0 \\
K(t, x, x) \Lambda(t, x)-\Lambda(t, x) K(t, x, x)=-M^{1}(t, x)
\end{array} \tag{39}
\end{array}\right.
$$

(ii) $G^{2}$ is defined by

$$
G^{2}=\left(\begin{array}{ll}
G_{--}^{2} & 0 \\
G_{+-}^{2} & 0
\end{array}\right)
$$

with

$$
\left\{\begin{array}{l}
G_{--}^{2}(t, x)=-K_{--}(t, x, 0) \Lambda_{--}(t, 0)-K_{-+}(t, x, 0) \Lambda_{++}(t, 0) Q^{1}(t)  \tag{40}\\
G_{+-}^{2}(t, x)=-K_{+-}(t, x, 0) \Lambda_{--}(t, 0)-K_{++}(t, x, 0) \Lambda_{++}(t, 0) Q^{1}(t)
\end{array}\right.
$$

Finally, the stability property (35) is clearly satisfied since, at every fixed $t \geq 0$, the Volterra transformation (37) defines an injective map of $L^{2}(0,1)^{n}$ (see e.g. Hoc73, Theorem 2.6]).

### 2.2.2 The kernel equations

We can prove that there exists $K \in C^{0}(\overline{\mathcal{T}})^{n \times n} \cap L^{\infty}(\mathcal{T})^{n \times n}$ that satisfies the so-called "kernel equations" (39) in the sense of broad solutions. However, it is in general not enough to deduce stability results for the initial system $(M, 0, F, Q)$ since the investigation of the stability properties of the system $\left(0, G^{2}, F^{2}, Q^{1}\right)$ is not an easier task without knowing any more information about it. The breakthrough idea of the conference paper HDM15 in the time-independent case (see also HDMVK16, HVDMK19]) was to construct a solution $K$ to the kernel equations which, in addition, yields a simpler structure for the matrix $G_{--}^{2}$ defined in $(40)$. This is the key point to prove stability results for the system $\left(0, G^{2}, F^{2}, Q^{1}\right)$ (see Remark 2.5 and Section 2.3 below). Such a construction is possible by adding some conditions for $K_{--}$at $(t, x, 0)$ (see 40$)$ but the price to pay is that it introduces discontinuities for $K_{--}$, so that $K$ will not be globally $C^{0}$ anymore but only piecewise $C^{0}$ in general. We will prove the following result:

Theorem 2.6. There exists a $n \times n$ matrix-valued function $K=\left(k_{i j}\right)_{1 \leq i, j \leq n}$ such that:
(i) $K \in L^{\infty}(\mathcal{T})^{n \times n}$.
(ii) For every $i, j \in\{1, \ldots, n\}$ with $j \notin\{i+1, \ldots, m\}$, we have $k_{i j} \in C^{0}(\overline{\mathcal{T}})$.
(iii) For every $i, j \in\{1, \ldots, m\}$ with $i<j$, we have $k_{i j} \in C^{0}\left(\overline{\mathcal{T}_{i j}^{-}}\right) \cap C^{0}\left(\overline{\mathcal{T}_{i j}^{+}}\right)$, where (see Figure 1 )

$$
\begin{aligned}
& \mathcal{T}_{i j}^{-}=\left\{(t, x, \xi) \in \mathcal{T}, \quad \xi<\psi_{i j}(t, x)\right\} \\
& \mathcal{T}_{i j}^{+}=\left\{(t, x, \xi) \in \mathcal{T}, \quad \xi>\psi_{i j}(t, x)\right\}
\end{aligned}
$$

where $\psi_{i j} \in C^{1}([0,+\infty) \times[0,1])$ satisfies the following semi-linear hyperbolic equation for every
v) $K$ is a broad solution of (39) in $\mathcal{T}$ (the exact meaning of this statement will be detailed during the proof of the theorem, in Section (3.2 below).
(v) For every $t \geq 0$ and $x \in[0,1]$, the matrix $G_{--}^{2}(t, x)$ defined in (40) is strictly lower triangular, i.e. it satisfies (34) (it then follows from (ii) that $\left.G_{-\_}^{2} \in C^{0}([0,+\infty) \times[0,1])^{m \times m}\right)$.

The proof of Theorem 2.6 is one of the main technical difficulties of this article and it is postponed to Section 3.2 below for the sake of the presentation. We conclude this section with some important remarks.

Remark 2.7. Let us rewrite the second condition of (39) component-wise:

$$
\begin{equation*}
\left(\lambda_{j}(t, x)-\lambda_{i}(t, x)\right) k_{i j}(t, x, x)=-m_{i j}^{1}(t, x) . \tag{42}
\end{equation*}
$$

Therefore, we see that for $i=j$ we shall necessarily have $m_{i i}^{1}=0$ and it explains why we had to perform a preliminary transformation in Section 2.1 to remove these terms (otherwise the equation (42), and thus the kernel equations (39), have no solution).

Remark 2.8. It is in general not possible to solve 3 with $G_{--}^{2}=0$, unless $m=1$.
Remark 2.9. Observe that, with the regularity stated in Theorem 2.6 we have in particular that, for every $w \in C^{0}\left(\left[t^{0},+\infty\right) ; L^{2}(0,1)^{n}\right), t^{0} \geq 0$,

$$
(t, x) \mapsto \int_{0}^{x} K(t, x, \xi) w(t, \xi) d \xi \in C^{0}\left(\left[t^{0},+\infty\right) \times[0,1]\right)^{n}
$$

This follows from Lebesgue's dominated convergence theorem. This shows that $\gamma$ defined by (37) has the good regularity to be a broad solution (see Definition A.1), if so has $w$.
Remark 2.10. Observe that the condition (iii) shows that the kernel has possible discontinuities on $\xi=\psi_{i j}(t, x)$ for $i<j \leq m$. Besides, these discontinuities also depend on the component of the kernel that we consider. The appearance of such discontinuities is explained by the requirement of the last condition (v) because we somehow force two boundary conditions at the points ( $t, 0,0$ ), one by the condition already required in (39) (which concerns $i \neq j$, see Remark 2.7) and another one by (v) (which only concerns $i \leq j \leq m$ ). This results in discontinuities along the characteristics passing through these points. Note that this also complicates the justification of formal computations that we performed above since regularity problems will occur during the computation of the following term (when $i<j \leq m$ ):

$$
\frac{\partial}{\partial t}\left(\int_{0}^{x} k_{i j}(t, x, \xi) w_{j}(t, \xi) d \xi\right)+\lambda_{i}(t, x) \frac{\partial}{\partial x}\left(\int_{0}^{x} k_{i j}(t, x, \xi) w_{j}(t, \xi) d \xi\right)
$$

More precisely, writing $\int_{0}^{x}=\int_{0}^{\psi_{i j}(t, x)}+\int_{\psi_{i j}(t, x)}^{x}$ and using integration by parts, we see that the following jump terms notably appear:

$$
\begin{aligned}
&\left(\frac{\partial \psi_{i j}}{\partial t}(t, x)-\lambda_{j}\left(t, \psi_{i j}(t, x)\right)+\lambda_{i}(t, x) \frac{\partial \psi_{i j}}{\partial x}(t, x)\right) \\
& \times\left(k_{i j}^{-}\left(t, x, \psi_{i j}(t, x)\right)-k_{i j}^{+}\left(t, x, \psi_{i j}(t, x)\right)\right) w_{j}\left(t, \psi_{i j}(t, x)\right)
\end{aligned}
$$

where $k_{i j}^{-}$(resp. $k_{i j}^{+}$) denotes the trace of the restriction of $k_{i j}$ to $\partial \mathcal{T}_{i j}^{-}$(resp $\partial \mathcal{T}_{i j}^{+}$). This is why it is crucial to precise that $\psi_{i j}$ solves the first equation in (41) so that such undesired terms vanish in the end. In the case of time-independent systems, we have in fact

$$
\psi_{i j}(t, x)=\phi_{j}^{-1}\left(\phi_{i}(x)\right),
$$

where we introduced $\phi_{\ell}(x)=\int_{0}^{x} \frac{1}{-\lambda_{\ell}(\xi)} d \xi$ for $\ell \in\{1, \ldots, m\}\left(\psi_{i j}\right.$ is well defined because $\left.i<j\right)$. This is the same function as in HVDMK19, (A.1)].


Figure 1: 2D cross-section of the domain $\mathcal{T}$ at a fixed $t$

### 2.3 Fredholm integral transformation

We recall that at the moment we already know that the system $\left(0, G^{2}, F^{2}, Q^{1}\right)$ of Proposition 2.4 is finite-time stable if we take $F^{2}=0$, but only with a settling time which is strictly larger than $T_{\text {unif }}(\Lambda)$ (unless $m=1$ ), see Remark 2.5. In this section, we perform a third and last transformation to remove the coupling term $G_{--}^{2}$ in the system $\left(0, G^{2}, F^{2}, Q^{1}\right)$ and we show that the resulting system has the desired stability properties. More precisely, the goal of this section is to establish the two following results:

Proposition 2.11. There exists $F^{2} \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ such that the following property holds for every $T>0$ :

$$
\begin{align*}
&\left(0, G^{3}, 0, Q^{1}\right) \text { is finite-time stable with settling time } T \\
& \Longrightarrow \quad\left(0, G^{2}, F^{2}, Q^{1}\right) \text { is finite-time stable with settling time } T \tag{43}
\end{align*}
$$

where

$$
G^{3}=\left(\begin{array}{cc}
0 & 0  \tag{44}\\
G_{+-}^{2} & 0
\end{array}\right)
$$

Proposition 2.12. The system $\left(0, G^{3}, 0, Q^{1}\right)$ is finite-time stable with settling time $T_{\text {unif }}(\Lambda)$ defined by (22).

Note that the proof of our main result - Theorem 1.5 - will then be complete (recall Propositions 2.1 and 2.4, except for the $\tau$-periodicity statement which will be studied later on in Section 3.3 .

### 2.3.1 Finite-time stability of the system $\left(0, G^{3}, 0, Q^{1}\right)$

In this section we prove Proposition 2.12 in four steps.

1) Let $t^{0} \geq 0$ be fixed. From the very definition of broad solution (see Definition A.1 and the simple structure of $G^{3}$, we see that the first $m$ components of the system vanish at time $t^{0}+T_{\text {unif }}(\Lambda)$ if (recall that the feedback is equal to zero)

$$
\begin{equation*}
s_{i}^{\mathrm{in}}\left(t^{0}+T_{\mathrm{unif}}(\Lambda), x\right)>t^{0}, \quad \forall x \in[0,1], \quad \forall i \in\{1, \ldots, m\} \tag{45}
\end{equation*}
$$

and the remaining $p$ components of the system vanish at time $t^{0}+T_{\text {unif }}(\Lambda)$ if

$$
\left\{\begin{array}{l}
s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), x\right)>t^{0}, \quad \forall x \in[0,1], \quad \forall i \in\{m+1, \ldots, n\}  \tag{46}\\
s_{j}^{\text {in }}\left(s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), x\right), 0\right)>t^{0}, \quad \forall x \in[0,1), \quad \forall j \in\{1, \ldots, m\}, \quad \forall i \in\{m+1, \ldots, n\}
\end{array}\right.
$$

2) First of all, observe that, from (20), 17) and (16) we have the following inverse formula for every $t, \bar{t} \in \mathbb{R}$ :

$$
\left\{\begin{array}{rll}
s_{i}^{\text {in }}(t, 0)>\bar{t} & \Longleftrightarrow t>s_{i}^{\text {out }}(\bar{t}, 1), & \text { if } i \in\{1, \ldots, m\}  \tag{47}\\
s_{i}^{\text {in }}(t, 1) \geq \bar{t} & \Longleftrightarrow t \geq s_{i}^{\text {out }}(\bar{t}, 0), & \text { if } i \in\{m+1, \ldots, n\}
\end{array}\right.
$$

3) Let us establish (45). Let then $i \in\{1, \ldots, m\}$ be fixed. We have:

$$
\begin{aligned}
s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), x\right)>t^{0}, \quad \forall x \in[0,1] & \Longleftrightarrow s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), 0\right)>t^{0}, \quad(\text { by } \\
& \Longleftrightarrow t^{0}+T_{\text {unif }}(\Lambda)>s_{i}^{\text {out }}\left(t^{0}, 1\right), \quad(\text { by }),
\end{aligned}
$$

and this last statement holds true since, by definition of $T_{\text {unif }}(\Lambda)$ and $\sqrt{15}-16$, we have, for an arbitrary $j \in\{m+1, \ldots, n\}$,

$$
t^{0}+T_{\text {unif }}(\Lambda) \geq s_{j}^{\text {out }}\left(s_{i}^{\text {out }}\left(t^{0}, 1\right), 0\right)>s_{i}^{\text {out }}\left(t^{0}, 1\right)
$$

4) Let us now establish (46). We focus on the second inequality since the first one is obtained similarly to 45 . Let then $i \in\{m+1, \ldots, n\}$ and $j \in\{1, \ldots, m\}$ be fixed. We have:

$$
\begin{aligned}
s_{j}^{\text {in }}\left(s _ { i } ^ { \text { in } } \left(t^{0}+T_{\text {unif }}(\Lambda),\right.\right. & x), 0)>t^{0}, \quad \forall x \in[0,1) \\
& \Longleftrightarrow s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), x\right)>s_{j}^{\text {out }}\left(t^{0}, 1\right), \quad \forall x \in[0,1), \quad(\text { by }(47)), \\
& \Longleftrightarrow s_{i}^{\text {in }}\left(t^{0}+T_{\text {unif }}(\Lambda), 1\right) \geq s_{j}^{\text {out }}\left(t^{0}, 1\right), \quad(\text { by }(20)), \\
& \Longleftrightarrow t^{0}+T_{\text {unif }}(\Lambda) \geq s_{i}^{\text {out }}\left(s_{j}^{\text {out }}\left(t^{0}, 1\right), 0\right), \quad(\text { by (47) }),
\end{aligned}
$$

and this last statement holds true by definition of $T_{\text {unif }}(\Lambda)$.

### 2.3.2 Proof of Proposition 2.11

We start the proof with some computations. We will show that we can transform a solution of $\left(0, G^{3}, 0, Q^{1}\right)$ into a solution of $\left(0, G^{2}, F^{2}, Q^{1}\right)$ (note the difference in the order of the transformation with respect to the previous sections and see Remark 2.15 below for the reason). Let then $z$ be the solution to the system $\left(0, G^{3}, 0, Q^{1}\right)$ with initial data $z^{0}$. Inspired by the work CHO17 mentioned before (for time-independent systems), we propose to use a Fredholm integral transformation as follows:

$$
\begin{equation*}
\gamma(t, x)=z(t, x)-\int_{0}^{1} H(t, x, \xi) z(t, \xi) d \xi \tag{48}
\end{equation*}
$$

where we suppose for the moment that the kernel $H$ is smooth on $\overline{\mathcal{R}}$, where $\mathcal{R}$ is the infinite rectangular prism defined by

$$
\mathcal{R}=(0,+\infty) \times(0,1) \times(0,1)
$$

Let us now perform some formal computations and see what $\gamma$ can solve. We have

$$
\begin{aligned}
\frac{\partial \gamma}{\partial t}(t, x)+ & \Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=\frac{\partial z}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial z}{\partial x}(t, x) \\
& -\int_{0}^{1} \frac{\partial H}{\partial t}(t, x, \xi) z(t, \xi) d \xi-\int_{0}^{1} H(t, x, \xi) \frac{\partial z}{\partial t}(t, \xi) d \xi-\Lambda(t, x) \int_{0}^{1} \frac{\partial H}{\partial x}(t, x, \xi) z(t, \xi) d \xi
\end{aligned}
$$

Using the equation satisfied by $z$, we obtain

$$
\begin{aligned}
\frac{\partial \gamma}{\partial t}(t, x) & +\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=G^{3}(t, x) z(t, 0)-\int_{0}^{1} \frac{\partial H}{\partial t}(t, x, \xi) z(t, \xi) d \xi \\
& -\int_{0}^{1} H(t, x, \xi)\left(-\Lambda(t, \xi) \frac{\partial z}{\partial \xi}(t, \xi)+G^{3}(t, \xi) z(t, 0)\right) d \xi-\Lambda(t, x) \int_{0}^{1} \frac{\partial H}{\partial x}(t, x, \xi) z(t, \xi) d \xi
\end{aligned}
$$

Integrating by parts the third term of the right hand side, we obtain

$$
\begin{aligned}
\frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=\int_{0}^{1} & \left(-\frac{\partial H}{\partial t}(t, x, \xi)-\frac{\partial H}{\partial \xi}(t, x, \xi) \Lambda(t, \xi)-H(t, x, \xi) \frac{\partial \Lambda}{\partial \xi}(t, \xi)\right. \\
& \left.-\Lambda(t, x) \frac{\partial H}{\partial x}(t, x, \xi)\right) z(t, \xi) d \xi \\
+H(t, x, 1) \Lambda(t, 1) z(t, 1)+ & \left(G^{3}(t, x)-H(t, x, 0) \Lambda(t, 0)-\int_{0}^{1} H(t, x, \xi) G^{3}(t, \xi) d \xi\right) z(t, 0)
\end{aligned}
$$

Using the formula 48 with $x=0$ we finally obtain

$$
\begin{aligned}
& \frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=\int_{0}^{1}\left(-\frac{\partial H}{\partial t}(t, x, \xi)-\frac{\partial H}{\partial \xi}(t, x, \xi) \Lambda(t, \xi)-H(t, x, \xi) \frac{\partial \Lambda}{\partial \xi}(t, \xi)\right. \\
& \left.-\Lambda(t, x) \frac{\partial H}{\partial x}(t, x, \xi)+\left(G^{3}(t, x)-H(t, x, 0) \Lambda(t, 0)-\int_{0}^{1} H(t, x, \zeta) G^{3}(t, \zeta) d \zeta\right) H(t, 0, \xi)\right) z(t, \xi) d \xi \\
& \quad+H(t, x, 1) \Lambda(t, 1) z(t, 1)+\left(G^{3}(t, x)-H(t, x, 0) \Lambda(t, 0)-\int_{0}^{1} H(t, x, \xi) G^{3}(t, \xi) d \xi\right) \gamma(t, 0)
\end{aligned}
$$

Since

$$
\begin{equation*}
z_{-}(t, 1)=0 \tag{49}
\end{equation*}
$$

the boundary term $H(t, x, 1) \Lambda(t, 1) z(t, 1)$ vanishes if we require that $H$ satisfies

$$
H_{-+}(t, x, 1)=H_{++}(t, x, 1)=0
$$

On the other hand, $\gamma$ and $z$ satisfy the same boundary condition at $x=0$ provided that

$$
H(t, 0, \xi)=0
$$

Finally, at $x=1$, we have (recall 49 )

$$
\gamma_{-}(t, 1)-\int_{0}^{1} F^{2}(t, \xi) \gamma(t, \xi) d \xi=\int_{0}^{1}\left(-H_{-}(t, 1, \xi)-F^{2}(t, \xi)+\int_{0}^{1} F^{2}(t, \zeta) H(t, \zeta, \xi) d \zeta\right) z(t, \xi) d \xi
$$

where $H_{-}$denotes again the $m \times n$ sub-matrix of $H$ formed by its first $m$ rows. Thus, we see that $\gamma$ satisfies at $x=1$ the boundary condition $\gamma_{-}(t, 1)=\int_{0}^{1} F^{2}(t, \xi) \gamma(t, \xi) d \xi$ if $F^{2}(t, \cdot)$ satisfies the following Fredholm integral equation (at $t$ fixed):

$$
\begin{equation*}
F^{2}(t, \xi)-\int_{0}^{1} F^{2}(t, \zeta) H(t, \zeta, \xi) d \zeta=-H_{-}(t, 1, \xi) \tag{50}
\end{equation*}
$$

In summary, $\gamma$ defined by (48) is the solution of $\left(0, G^{2}, F^{2}, Q^{1}\right)$ with state-feedback gain function $F^{2}$ satisfying (50) (whenever it exists) and initial data $\gamma^{0}(x)=z^{0}(x)-\int_{0}^{1} H\left(t^{0}, x, \xi\right) z^{0}(\xi) d \xi$ if we have the following two properties:
(i) For every $(t, x, \xi) \in \mathcal{R}$,

$$
\left\{\begin{array}{l}
\frac{\partial H}{\partial t}(t, x, \xi)+\Lambda(t, x) \frac{\partial H}{\partial x}(t, x, \xi)+\frac{\partial H}{\partial \xi}(t, x, \xi) \Lambda(t, \xi)+H(t, x, \xi) \frac{\partial \Lambda}{\partial \xi}(t, \xi)=0  \tag{51}\\
H_{-+}(t, x, 1)=H_{++}(t, x, 1)=H(t, 0, \xi)=0
\end{array}\right.
$$

(ii) $G^{3}$ satisfies the Fredholm integral equation

$$
\begin{equation*}
G^{3}(t, x)-\int_{0}^{1} H(t, x, \xi) G^{3}(t, \xi) d \xi=G^{2}(t, x)+H(t, x, 0) \Lambda(t, 0) \tag{52}
\end{equation*}
$$

Finally, the stability property $(43)$ is clearly satisfied if, for every $t \geq 0$, the Fredholm transformation (48) defines a surjective map of $L^{2}(0,1)^{n}$.

It remains to prove the existence of $F^{2}$ and $H$ satisfying the above properties and so that the Fredholm transformation (48) is invertible (let us recall that, unlike Volterra transformations of the second kind, Fredholm transformations are not always invertible). Note that $H=0$ is a solution of (51). Taking into account the very particular structure (36) of $G^{2}$, this motivates our attempt to look for a kernel $H$ with the following simple structure:

$$
H=\left(\begin{array}{cc}
H_{--} & 0  \tag{53}\\
0 & 0
\end{array}\right)
$$

This structure implies that the Fredholm equation (52) is equivalent to

$$
\left\{\begin{array}{l}
G_{--}^{3}(t, x)-\int_{0}^{1} H_{--}(t, x, \xi) G_{--}^{3}(t, \xi) d \xi=G_{--}^{2}(t, x)+H_{--}(t, x, 0) \Lambda_{--}(t, 0) \\
G_{-+}^{3}(t, x)-\int_{0}^{1} H_{--}(t, x, \xi) G_{-+}^{3}(t, \xi) d \xi=0 \\
G_{+-}^{3}(t, x)=G_{+-}^{2}(t, x) \\
G_{++}^{3}(t, x)=0
\end{array}\right.
$$

These equations are easily solved by taking $G_{--}^{3}=G_{-+}^{3}=0$ (so that $G^{3}$ is indeed given by (44)) if we impose the following condition for $H_{--}$at $(t, x, 0)$ :

$$
H_{--}(t, x, 0)=-G_{--}^{2}(t, x) \Lambda_{--}(t, 0)^{-1}
$$

We point out that this last condition may introduce discontinuities in the kernel, because of possible compatibility conditions at $(t, 0,0)$ with the previous requirement that $H_{--}(t, 0, \xi)=0$.

Finally, since $G_{--}^{2}$ is in fact strictly lower triangular (34), we also look for $H_{--}$with the same structure. Note that this structure in particular ensures that the Fredholm transformation (48) is invertible and that the Fredholm equation (50) always has a unique solution $F^{2} \in L^{\infty}((0,+\infty) \times$ $(0,1))^{m \times n}$, provided that $H_{--} \in L^{\infty}(\mathcal{R})^{m \times m}$ (see for instance [CHO17, Appendix] for more details). This property was a priori not guaranteed without additional information (we emphasize again that Fredholm transformations are not always invertible).

The final step is to prove the existence of $H_{--}$that satisfies all the properties mentioned above. This is the goal of the following theorem, the proof of which is given in Section 3.1 below.
Theorem 2.13. There exists a $m \times m$ matrix-valued function $H_{--}=\left(h_{i j}\right)_{1 \leq i, j \leq m}$ such that:
(i) $H_{--} \in L^{\infty}(\mathcal{R})^{m \times m}$.
(ii) For $i \leq j$, we have $h_{i j}=0$ (i.e. $H_{--}$is strictly lower triangular).
(iii) For $i>j$, we have $h_{i j} \in C^{0}\left(\overline{\mathcal{R}_{i j}^{-}}\right) \cap C^{0}\left(\overline{\mathcal{R}_{i j}^{+}}\right)$, where

$$
\begin{aligned}
& \mathcal{R}_{i j}^{-}=\left\{(t, x, \xi) \in \mathcal{R}, \quad \xi<\psi_{i j}(t, x)\right\}, \\
& \mathcal{R}_{i j}^{+}=\left\{(t, x, \xi) \in \mathcal{R}, \quad \xi>\psi_{i j}(t, x)\right\},
\end{aligned}
$$

where $\psi_{i j} \in C^{1}([0,+\infty) \times[0,1])$ satisfies the semi-linear hyperbolic equation 41.
(iv) $H_{--}$is the unique broad solution in $\mathcal{R}$ of the system

$$
\left\{\begin{array}{l}
\frac{\partial H_{--}}{\partial t}(t, x, \xi)+\Lambda_{--}(t, x) \frac{\partial H_{--}}{\partial x}(t, x, \xi)+\frac{\partial H_{--}}{\partial \xi}(t, x, \xi) \Lambda_{--}(t, \xi)  \tag{54}\\
\\
\\
\\
\\
H_{--}(t, 0, \xi)=0 \\
H_{--}(t, x, 0)=-H_{--}(t, x, \xi) \frac{\partial \Lambda_{--}}{\partial \xi}(t, \xi)=0
\end{array}\right.
$$

(once again, the exact meaning of this statement will be detailed during the proof of the theorem, in Section 3.1 below).

This concludes the proof of Proposition 2.11

Remark 2.14. Observe once again that the kernel is discontinuous. This introduces some additional boundary terms along these discontinuities in the formal computations performed above but, as mentioned before in Remark 2.10, these terms cancel each other out thanks to the equation satisfied by $\psi_{i j}$ in (41). As in the time-independent case ([CHO17, Section 3]), the system (54) is easy to solve and its solution is even explicit (see (60)-58) below). sense).

Writing (54) component-wise, this gives

$$
\left\{\begin{array}{l}
\frac{\partial h_{i j}}{\partial t}(t, x, \xi)+\lambda_{i}(t, x) \frac{\partial h_{i j}}{\partial x}(t, x, \xi)+\lambda_{j}(t, \xi) \frac{\partial h_{i j}}{\partial \xi}(t, x, \xi)+\frac{\partial \lambda_{j}}{\partial \xi}(t, \xi) h_{i j}(t, x, \xi)=0  \tag{56}\\
h_{i j}(t, 0, \xi)=0 \\
h_{i j}(t, x, 0)=-\frac{g_{i j}^{2}(t, x)}{\lambda_{j}(t, 0)}
\end{array}\right.
$$

Since we see that the equations are uncoupled, we can fix the indices $i, j$ for the remainder of Section 3.1

$$
i, j \in\{1, \ldots, m\} \text { are fixed. }
$$

### 3.1.1 The characteristics of (56)

For each $(t, x, \xi) \in \mathbb{R}^{3}$ fixed, we introduce the characteristic curve $\chi_{i j}(\cdot ; t, x)$ associated with the hyperbolic equation (56) passing through the point $(t, x, \xi)$, i.e.

$$
\chi_{i j}(s ; t, x, \xi)=\left(s, \chi_{i}(s ; t, x), \chi_{j}(s ; t, \xi)\right), \quad \forall s \in \mathbb{R}
$$

where we recall that $\chi_{i}$ and $\chi_{j}$ are defined in (9). For every $(t, x, \xi) \in \mathcal{R}$, we have

$$
\chi_{i j}(s ; t, x, \xi) \in \mathcal{R}, \quad \forall s \in\left(s_{i j}^{\text {in }}(t, x, \xi), s_{i j}^{\text {out }}(t, x, \xi)\right)
$$

where we introduced

$$
s_{i j}^{\text {in }}(t, x, \xi)=\max \left\{0, s_{i}^{\text {in }}(t, x), s_{j}^{\text {in }}(t, \xi)\right\}>0, \quad s_{i j}^{\text {out }}(t, x, \xi)=\min \left\{s_{i}^{\text {out }}(t, x), s_{j}^{\text {out }}(t, \xi)\right\}
$$

Since the speeds $\lambda_{i}, \lambda_{j}$ are negative $(i, j \leq m)$, when $s$ is increasing, $s \mapsto \chi_{i}(s ; t, x), s \mapsto \chi_{j}(s ; t, \xi)$ are decreasing. Therefore, the associated characteristic $\chi_{i j}(\cdot ; t, x, \xi)$ will exit the domain $\mathcal{R}$ through the planes $x=0$ or $\xi=0$. This is why we can impose boundary conditions at $(t, 0, \xi)$ and $(t, x, 0)$ (see 56 ) and this is why it is enough to (uniquely) determine a solution on $\mathcal{R}$. To be more precise, we can split $\mathcal{R}$ into three disjoint subsets:

$$
\mathcal{R}=\mathcal{R}_{i j}^{+} \cup \mathcal{R}_{i j}^{-} \cup \mathcal{D}_{i j}
$$

where

$$
\begin{array}{ll}
\mathcal{R}_{i j}^{+}=\{(t, x, \xi) \in \mathcal{R}, & \left.s_{i}^{\text {out }}(t, x)<s_{j}^{\text {out }}(t, \xi)\right\} \\
\mathcal{R}_{i j}^{-}=\{(t, x, \xi) \in \mathcal{R}, & \left.s_{i}^{\text {out }}(t, x)>s_{j}^{\text {out }}(t, \xi)\right\} \\
\mathcal{D}_{i j}=\{(t, x, \xi) \in \mathcal{R}, & \left.s_{i}^{\text {out }}(t, x)=s_{j}^{\text {out }}(t, \xi)\right\}
\end{array}
$$

With these notations, the characteristic $\chi_{i j}(\cdot ; t, x, \xi)$ will either exit the domain $\mathcal{R}$ through the plane $x=0$ if $(t, x, \xi) \in \mathcal{R}_{i j}^{+}$or through the plane $\xi=0$ if $(t, x, \xi) \in \mathcal{R}_{i j}^{-}$:

## Proposition 3.1.

(i) For every $(t, x, \xi) \in \mathcal{R}_{i j}^{+}$, we have $\chi_{i j}(s ; t, x, \xi) \in \mathcal{R}_{i j}^{+}$for every $s \in\left(t, s_{i}^{\text {out }}(t, x)\right)$.
(ii) For every $(t, x, \xi) \in \mathcal{R}_{i j}^{-}$, we have $\chi_{i j}(s ; t, x, \xi) \in \mathcal{R}_{i j}^{-}$for every $s \in\left(t, s_{j}^{\text {out }}(t, \xi)\right)$.
(iii) For every $(t, x, \xi) \in \mathcal{D}_{i j}$, we have $\chi_{i j}(s ; t, x, \xi) \in \mathcal{D}_{i j}$ for every $s \in\left(t, s_{i}^{\text {out }}(t, x)\right)=\left(t, s_{j}^{\text {out }}(t, \xi)\right)$.

These three points directly follow from (17).

### 3.1.2 Existence and regularity of a solution to 5

Writing the solution of (56) along the characteristic curve $\chi_{i j}(s ; t, x, \xi)$ for $s \in\left[s_{i j}^{\text {in }}(t, x, \xi), s_{i j}^{\text {out }}(t, x, \xi)\right]$ and using the boundary conditions, we obtain the following ODE:

$$
\left\{\begin{array}{l}
\frac{d}{d s} h_{i j}\left(\chi_{i j}(s ; t, x, \xi)\right)=-\frac{\partial \lambda_{j}}{\partial \xi}\left(s, \chi_{j}(s ; t, \xi)\right) h_{i j}\left(\chi_{i j}(s ; t, x, \xi)\right)  \tag{57}\\
h_{i j}\left(\chi_{i j}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x, \xi\right)\right)=b_{i j}(t, x, \xi)
\end{array}\right.
$$

where

$$
b_{i j}(t, x, \xi)=\left\{\begin{array}{cc}
0 & \text { if }(t, x, \xi) \in \mathcal{R}_{i j}^{+}  \tag{58}\\
-\frac{g_{i j}^{2}\left(s_{j}^{\text {out }}(t, \xi), \chi_{i}\left(s_{j}^{\mathrm{out}}(t, \xi) ; t, x\right)\right)}{\lambda_{j}\left(s_{j}^{\text {out }}(t, \xi), 0\right)} & \text { if }(t, x, \xi) \in \mathcal{R}_{i j}^{-}
\end{array}\right.
$$

Integrating this ODE over $\left[t, s_{i j}^{\text {out }}(t, x, \xi)\right]$ yields the integral equation

$$
\begin{equation*}
h_{i j}(t, x, \xi)=b_{i j}(t, x, \xi)+\int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} \frac{\partial \lambda_{j}}{\partial \xi}\left(s, \chi_{j}(s ; t, \xi)\right) h_{i j}\left(\chi_{i j}(s ; t, x, \xi)\right) d s \tag{59}
\end{equation*}
$$

In this case, this integral equation is very easily solved by taking (as it is in fact directly seen from the ODE (57)):

$$
\begin{equation*}
h_{i j}(t, x, \xi)=b_{i j}(t, x, \xi) e^{\int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} \frac{\partial \lambda_{j}}{\partial \xi}\left(s, \chi_{j}(s ; t, \xi)\right) d s} . \tag{60}
\end{equation*}
$$

Clearly, $h_{i j}=0$ for $i \leq j$ (i.e. $H_{--}$is indeed strictly lower triangular) since $g_{i j}^{2}=0$ for such indices (see (34). Obviously, $h_{i j} \in C^{0}\left(\overline{\mathcal{R}_{i j}^{+}}\right) \cap L^{\infty}\left(\mathcal{R}_{i j}^{+}\right)$. On the other hand, thanks in particular to the regularities 41,418 , the bounds 19$)$ and the assumption $\frac{\partial \Lambda}{\partial x} \in L^{\infty}((0,+\infty) \times(0,1))^{n \times n}$, we can check that

$$
h_{i j} \in C^{0}\left(\overline{\mathcal{R}_{i j}^{-}}\right) \cap L^{\infty}\left(\mathcal{R}_{i j}^{-}\right) .
$$

### 3.1.3 Characterization of $\mathcal{R}_{i j}^{ \pm}$and $\mathcal{D}_{i j}$

Let us now show that

$$
\begin{equation*}
\mathcal{R}_{i j}^{-}=\left\{(t, x, \xi) \in \mathcal{R}, \quad \xi<\psi_{i j}(t, x)\right\}, \tag{61}
\end{equation*}
$$

$$
\mathcal{R}_{i j}^{+}=\left\{(t, x, \xi) \in \mathcal{R}, \quad \xi>\psi_{i j}(t, x)\right\}
$$

where $\psi_{i j} \in C^{1}([0,+\infty) \times[0,1])$ satisfies the semi-linear hyperbolic equation (41). First of all, it follows from $(20)$ and the implicit function theorem that there exists a function $\psi_{i j} \in C^{1}([0,+\infty) \times$ $[0,1]), 0 \leq \psi_{i j} \leq 1$, such that

$$
\begin{equation*}
s_{i}^{\text {out }}(t, x)=s_{j}^{\text {out }}(t, \xi) \quad \Longleftrightarrow \quad \xi=\psi_{i j}(t, x) \tag{62}
\end{equation*}
$$

This shows that

$$
\begin{equation*}
\mathcal{D}_{i j}=\left\{(t, x, \xi) \in \mathcal{R}, \quad \xi=\psi_{i j}(t, x)\right\} . \tag{63}
\end{equation*}
$$

On the other hand, thanks to $(20)$ and $(62)$ we have

$$
\xi>\psi_{i j}(t, x) \quad \Longleftrightarrow \quad s_{j}^{\text {out }}(t, \xi)>s_{j}^{\text {out }}\left(t, \psi_{i j}(t, x)\right)=s_{i}^{\text {out }}(t, x)
$$

This shows the equality (61) for $\mathcal{R}_{i j}^{+}$. The equality for $\mathcal{R}_{i j}^{-}$can be proved similarly.
It remains to show that $\psi_{i j}$ satisfies the semi-linear hyperbolic equation (41). This in fact follows from (63) and (iii) of Proposition 3.1. Indeed, thanks to these results, we have

$$
\chi_{j}\left(s ; t, \psi_{i j}(t, x)\right)=\psi_{i j}\left(s, \chi_{i}(s ; t, x)\right), \quad \forall s \in\left(t, s_{i}^{\text {out }}(t, x)\right)=\left(t, s_{j}^{\text {out }}\left(t, \psi_{i j}(t, x)\right)\right) .
$$

Taking the derivative of this identity at $s=t^{+}$, we immediately obtain the equation in (41). On the other hand, letting $s \rightarrow s_{i}^{\text {out }}(t, x)^{-}=s_{j}^{\text {out }}\left(t, \psi_{i j}(t, x)\right)^{-}$and then letting $x \rightarrow 0^{+}$, we obtain the second condition $\psi_{i j}(t, 0)=0$.

### 3.2 Kernel for the Volterra transformation

In this section we prove Theorem 2.6, that is we prove the existence of a suitably smooth matrixvalued function $K=\left(k_{i j}\right)_{1 \leq i, j \leq n}$ such that

$$
\left\{\begin{array}{l}
\frac{\partial K}{\partial t}(t, x, \xi)+\Lambda(t, x) \frac{\partial K}{\partial x}(t, x, \xi)+\frac{\partial K}{\partial \xi}(t, x, \xi) \Lambda(t, \xi)+K(t, x, \xi) \widetilde{M}^{1}(t, \xi)=0  \tag{64}\\
K(t, x, x) \Lambda(t, x)-\Lambda(t, x) K(t, x, x)=-M^{1}(t, x)
\end{array}\right.
$$

where we introduced the notation

$$
\begin{equation*}
\widetilde{M}^{1}(t, \xi)=\frac{\partial \Lambda}{\partial \xi}(t, \xi)+M^{1}(t, \xi) \tag{65}
\end{equation*}
$$

Note in particular that $\widetilde{M}^{1} \in L^{\infty}((0,+\infty) \times(0,1))^{n \times n}$ thanks to the assumption (5). Besides, noticing (40), we also want the matrix

$$
\begin{equation*}
-K_{--}(t, x, 0) \Lambda_{--}(t, 0)-K_{-+}(t, x, 0) \Lambda_{++}(t, 0) Q^{1}(t) \quad \text { to be strictly lower triangular. } \tag{66}
\end{equation*}
$$

### 3.2.1 Preliminaries

Let us rewrite (64) by block. It is equivalent to the following four sub-systems:

$$
\begin{aligned}
& \left\{\begin{aligned}
& \frac{\partial K_{--}}{\partial t}(t, x, \xi)+\Lambda_{--}(t, x) \frac{\partial K_{--}}{\partial x}(t, x, \xi)+\frac{\partial K_{--}}{\partial \xi}(t, x, \xi) \Lambda_{--}(t, \xi) \\
&+K_{--}(t, x, \xi) \widetilde{M}_{--}^{1}(t, \xi)+K_{-+}(t, x, \xi) \widetilde{M}_{+-}^{1}(t, \xi)=0, \\
& K_{--}(t, x, x) \Lambda_{--}(t, x)-\Lambda_{--}(t, x) K_{--}(t, x, x)=-M_{--}^{1}(t, x),
\end{aligned}\right. \\
& \left\{\begin{aligned}
& \frac{\partial K_{-+}}{\partial t}(t, x, \xi)+\Lambda_{--}(t, x) \frac{\partial K_{-+}}{\partial x}(t, x, \xi)+\frac{\partial K_{-+}}{\partial \xi}(t, x, \xi) \Lambda_{++}(t, \xi) \\
&+K_{--}(t, x, \xi) \widetilde{M}_{-+}^{1}(t, \xi)+K_{-+}(t, x, \xi) \widetilde{M}_{++}^{1}(t, \xi)=0, \\
& K_{-+}(t, x, x) \Lambda_{++}(t, x)-\Lambda_{--}(t, x) K_{-+}(t, x, x)=-M_{-+}^{1}(t, x),
\end{aligned}\right. \\
& \left\{\begin{aligned}
& \frac{\partial K_{+-}}{\partial t}(t, x, \xi)+\Lambda_{++}(t, x) \frac{\partial K_{+-}}{\partial x}(t, x, \xi)+\frac{\partial K_{+-}}{\partial \xi}(t, x, \xi) \Lambda_{--}(t, \xi) \\
&+K_{+-}(t, x, \xi) \widetilde{M}_{--}^{1}(t, \xi)+K_{++}(t, x, \xi) \widetilde{M}_{+-}^{1}(t, \xi)=0, \\
& K_{+-}(t, x, x) \Lambda_{--}(t, x)-\Lambda_{++}(t, x) K_{+-}(t, x, x)=-M_{+-}^{1}(t, x),
\end{aligned}\right. \\
& \left\{\begin{aligned}
& \frac{\partial K_{++}}{\partial t}(t, x, \xi)+\Lambda_{++}(t, x) \frac{\partial K_{++}}{\partial x}(t, x, \xi)+\frac{\partial K_{++}}{\partial \xi}(t, x, \xi) \Lambda_{++}(t, \xi) \\
&+K_{+-}(t, x, \xi) \widetilde{M}_{-+}^{1}(t, \xi)+K_{++}(t, x, \xi) \widetilde{M}_{++}^{1}(t, \xi)=0, \\
& K_{++}(t, x, x) \Lambda_{++}(t, x)-\Lambda_{++}(t, x) K_{++}(t, x, x)=-M_{++}^{1}(t, x),
\end{aligned}\right.
\end{aligned}
$$

Remark 3.2. We see that $K_{--}$is coupled only with $K_{-+}$and that $K_{+-}$is coupled only with $K_{++}$. Moreover, the systems satisfied by $\left(K_{--}, K_{-+}\right)$and by ( $K_{+-}, K_{++}$) are similar. Therefore, from now on we only focus on the system satisfied by $\left(K_{--}, K_{-+}\right)$(note that the extra condition (66) only concerns this system). In addition, because of the nature of the coupling terms inside the domain (namely, matrix multiplication by the right), we see that the entries from different rows are not coupled. Therefore, for the rest of Section 3.2 we assume that

$$
i \in\{1, \ldots, m\} \text { is fixed. }
$$

Let us now rewrite the equations for $K_{--}$and $K_{-+}$component-wise. For convenience, we introduce

$$
\begin{equation*}
r_{i j}(t, x)=\frac{-m_{i j}^{1}(t, x)}{\lambda_{j}(t, x)-\lambda_{i}(t, x)} \quad(j \neq i) \tag{67}
\end{equation*}
$$

Note that $r_{i j} \in C^{0}([0,+\infty) \times[0,1])$. Moreover, $r_{i j} \in L^{\infty}((0,+\infty) \times(0,1))$ thanks to (4). We have:

### 3.2.2 The characteristics of (68)-(69)

For each $(t, x, \xi) \in \mathbb{R}^{3}$ fixed, we still denote by $\chi_{i j}(\cdot ; t, x, \xi)$ the characteristic curve associated with the hyperbolic system (68)-69) passing through the point $(t, x, \xi)$, i.e.

$$
\chi_{i j}(s ; t, x, \xi)=\left(s, \chi_{i}(s ; t, x), \chi_{j}(s ; t, \xi)\right), \quad \forall s \in \mathbb{R}
$$

We now need to find for which parameters $s$ the characteristic $\chi_{i j}(s ; t, x, \xi)$ stays in the domain $\mathcal{P}$ when $(t, x, \xi) \in \mathcal{P}$. To this end, we introduce the following sets for $j \in\{1, \ldots, m\}$ :

$$
\begin{aligned}
& \mathcal{P}_{i j}^{\text {in },+}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {in }}(t, x)<s_{j}^{\text {in }}(t, \xi)\right\}, \\
& \mathcal{P}_{i j}^{\text {in },-}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {in }}(t, x)>s_{j}^{\text {in }}(t, \xi)\right\}, \\
& \mathcal{D}_{i j}^{\text {in }}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {in }}(t, x)=s_{j}^{\text {in }}(t, \xi)\right\},
\end{aligned}
$$

and

$$
\begin{aligned}
& \mathcal{P}_{i j}^{\text {out },+}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {out }}(t, x)<s_{j}^{\text {out }}(t, \xi)\right\}, \\
& \mathcal{P}_{i j}^{\text {out },-}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {out }}(t, x)>s_{j}^{\text {out }}(t, \xi)\right\}, \\
& \mathcal{D}_{i j}^{\text {out }}=\left\{(t, x, \xi) \in \mathcal{P}, \quad s_{i}^{\text {out }}(t, x)=s_{j}^{\text {out }}(t, \xi)\right\} .
\end{aligned}
$$

As in Section 3.1.3 we can show that $\mathcal{P}_{i j}^{\text {out, }+} \cap \mathcal{T}=\mathcal{T}_{i j}^{+}$and $\mathcal{P}_{i j}^{\text {out },-} \cap \mathcal{T}=\mathcal{T}_{i j}^{-}$(we recall that $\mathcal{T}_{i j}^{+}$and $\mathcal{T}_{i j}^{-}$are defined in the statement of Theorem 2.6.

The following proposition gives precise information about the exit of the characteristics from the domain $\mathcal{P}$ (the proof is postponed to Appendix $B$ for the sake of the presentation; we refer to Figures 2, 3, 4 and 5 for a clarification of the geometric situation at a fixed $t$ ):

## Proposition 3.3.

(i) For every $j \in\{1, \ldots, i-1\}$, there exists a unique $s_{i j}^{\mathrm{in}} \in C^{0}(\overline{\mathcal{P}})$ with $(t, x, \xi) \mapsto t-s_{i j}^{\mathrm{in}}(t, x, \xi) \in$ $L^{\infty}(\mathcal{P})$ such that, for every $t \in \mathbb{R}$ and $0 \leq \xi<x<1$, we have $s_{i j}^{\mathrm{in}}(t, x, \xi)<t$ (and $s_{i j}^{\mathrm{in}}(t, x, \xi)=t$ otherwise) with

$$
\chi_{i j}(s ; t, x, \xi) \in \mathcal{P}, \quad \forall s \in\left(s_{i j}^{\operatorname{in}}(t, x, \xi), t\right)
$$

and

$$
\begin{cases}\chi_{j}\left(s_{i j}^{\mathrm{in}}(t, x, \xi) ; t, \xi\right)=\chi_{i}\left(s_{i j}^{\mathrm{in}}(t, x, \xi) ; t, x\right) & \text { if }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\mathrm{in},+}} \\ \chi_{i}\left(s_{i j}^{\mathrm{in}}(t, x, \xi) ; t, x\right)=1 & \text { if }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\mathrm{in},-}}\end{cases}
$$

(ii) For $j=i$, there exists a unique $s_{i i}^{\text {out }} \in C^{0}(\overline{\mathcal{P}})$ with $(t, x, \xi) \mapsto s_{i i}^{\text {out }}(t, x, \xi)-t \in L^{\infty}(\mathcal{P})$ such that, for every $t \in \mathbb{R}$ and $0<\xi \leq x \leq 1$, we have $s_{i i}^{\text {out }}(t, x, \xi)>t$ (and $s_{i i}^{\text {out }}(t, x, \xi)=t$ otherwise) and, if in addition $\xi<x$, then we have

$$
\chi_{i i}(s ; t, x, \xi) \in \mathcal{P}, \quad \forall s \in\left(t, s_{i i}^{\text {out }}(t, x, \xi)\right)
$$

and

$$
\chi_{i}\left(s_{i i}^{\text {out }}(t, x, \xi) ; t, \xi\right)=0
$$

(iii) For every $j \in\{i+1, \ldots, m\}$, there exists a unique $s_{i j}^{\text {out }} \in C^{0}(\overline{\mathcal{P}})$ with $(t, x, \xi) \mapsto s_{i j}^{\text {out }}(t, x, \xi)-$ $t \in L^{\infty}(\mathcal{P})$ such that, for every $t \in \mathbb{R}$ and $0<\xi<x \leq 1$, we have $s_{i j}^{\text {out }}(t, x, \xi)>t$ (and $s_{i j}^{\text {out }}(t, x, \xi)=t$ otherwise) with

$$
\chi_{i j}(s ; t, x, \xi) \in \mathcal{P}, \quad \forall s \in\left(t, s_{i j}^{\text {out }}(t, x, \xi)\right)
$$

and

$$
\begin{cases}\chi_{j}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, \xi\right)=\chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right) & \text { if }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out },+}} \\ \chi_{j}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, \xi\right)=0 & \text { if }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out },-}}\end{cases}
$$

(iv) For every $j \in\{m+1, \ldots, n\}$, there exists a unique $s_{i j}^{\text {out }} \in C^{0}(\overline{\mathcal{P}})$ with $(t, x, \xi) \mapsto s_{i j}^{\text {out }}(t, x, \xi)-$ $t \in L^{\infty}(\mathcal{P})$ such that, for every $t \in \mathbb{R}$ and $0 \leq \xi<x \leq 1$, we have $s_{i j}^{\text {out }}(t, x, \xi)>t$ (and $s_{i j}^{\text {out }}(t, x, \xi)=t$ otherwise) with

$$
\chi_{i j}(s ; t, x, \xi) \in \mathcal{P}, \quad \forall s \in\left(t, s_{i j}^{\text {out }}(t, x, \xi)\right)
$$

and

$$
\chi_{j}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, \xi\right)=\chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right)
$$



Figure 2: Definition of $s_{i j}^{\text {in }}$


Figure 3: Definition of $s_{i i}^{\text {out }}$
2) If $j=i$, then

$$
\left\{\begin{array}{l}
\frac{\partial k_{i i}}{\partial t}(t, x, \xi)+\lambda_{i}(t, x) \frac{\partial k_{i i}}{\partial x}(t, x, \xi)+\lambda_{i}(t, \xi) \frac{\partial k_{i i}}{\partial \xi}(t, x, \xi)+\sum_{\ell=1}^{n} k_{i \ell}(t, x, \xi) \widetilde{m}_{\ell i}^{1}(t, \xi)=0,  \tag{72}\\
k_{i i}(t, x, 0)=\sum_{\ell=1}^{p} k_{i, m+\ell}(t, x, 0) \widetilde{q}_{\ell i}^{1}(t) .
\end{array}\right.
$$

4) If $j \in\{m+1, \ldots, n\}$, then

$$
\left\{\begin{array}{l}
\frac{\partial k_{i j}}{\partial t}(t, x, \xi)+\lambda_{i}(t, x) \frac{\partial k_{i j}}{\partial x}(t, x, \xi)+\lambda_{j}(t, \xi) \frac{\partial k_{i j}}{\partial \xi}(t, x, \xi)+\sum_{\ell=1}^{n} k_{i \ell}(t, x, \xi) \widetilde{m}_{\ell j}^{1}(t, \xi)=0  \tag{74}\\
k_{i j}(t, x, x)=r_{i j}(t, x)
\end{array}\right.
$$

### 3.2.3 Transformation into integral equations

To prove the existence and uniqueness of the solution to the kernel equations $(71)-(74)$ on $\mathcal{P}$, we use the classical strategy that consists in transforming these hyperbolic equations into integral equations. Then, in the next subsection, we will prove that this system of integral equations has a unique solution by using a fixed-point argument and appropriate estimates.

Let us introduce

$$
\widetilde{k}_{i j}^{0}(t, x, \xi)= \begin{cases}r_{i j}\left(s_{i j}^{\text {in }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {in }}(t, x, \xi) ; t, x\right)\right) & \text { if } j \in\{1, \ldots, i-1\} \text { and }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {in, },+}}, \\ a_{i j}\left(s_{i j}^{\text {in }}(t, x, \xi), \chi_{j}\left(s_{i j}^{\text {in }}(t, x, \xi) ; t, \xi\right)\right) & \text { if } j \in\{1, \ldots, i-1\} \text { and }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {in, },}}, \\ r_{i j}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right)\right) & \text { if } j \in\{i+1, \ldots, m\} \text { and }(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out },+}}, \\ r_{i j}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right)\right) & \text { if } j \in\{m+1, \ldots, n\} .\end{cases}
$$

Thanks to the $C^{0}$-compatibility condition 70 , note in particular that

$$
\begin{equation*}
\widetilde{k}_{i j}^{0} \in C^{0}(\overline{\mathcal{P}}), \quad \forall j \in\{1, \ldots, i-1\} . \tag{75}
\end{equation*}
$$

Using now Proposition 3.3 we can obtain that

1) For $j \in\{1, \ldots, i-1\}$, integrating (71) along the characteristic curve $\chi_{i j}(s ; t, x, \xi)$ for $s \in$ $\left(s_{i j}^{\mathrm{in}}(t, x, \xi), t\right)$ yields the following integral equation:

$$
k_{i j}(t, x, \xi)=\widetilde{k}_{i j}^{0}(t, x, \xi)-\sum_{\ell=1}^{n} \int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s
$$

2) For $j=i$, integrating (72) along the characteristic curve $\chi_{i i}(s ; t, x, \xi)$ for $s \in\left(t, s_{i i}^{\text {out }}(t, x, \xi)\right)$ yields the following integral equation:

$$
\begin{align*}
& k_{i i}(t, x, \xi)=\sum_{\ell=1}^{p} k_{i, m+\ell}\left(s_{i i}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i i}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right) \widetilde{q}_{\ell i}^{1}\left(s_{i i}^{\text {out }}(t, x, \xi)\right) \\
&+\sum_{\ell=1}^{n} \int_{t}^{s_{i i}^{\text {out }}(t, x, \xi)} k_{i \ell}\left(\chi_{i i}(s ; t, x, \xi)\right) \widetilde{m}_{\ell i}^{1}\left(s, \chi_{i}(s ; t, \xi)\right) d s . \tag{76}
\end{align*}
$$

3) For $j \in\{i+1, \ldots, m\}$, integrating (73) along the characteristic curve $\chi_{i j}(s ; t, x, \xi)$ for $s \in$ $\left(t, s_{i j}^{\text {out }}(t, x, \xi)\right)$ yields the following integral equations:

$$
\begin{aligned}
& k_{i j}(t, x, \xi)=\widetilde{k}_{i j}^{0}(t, x, \xi) \\
&+\sum_{\ell=1}^{n} \int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s, \quad(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out },+}}
\end{aligned}
$$

and

$$
\begin{align*}
k_{i j}(t, x, \xi)= & \sum_{\ell=1}^{p} k_{i, m+\ell}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right) \widetilde{q}_{\ell j}^{1}\left(s_{i j}^{\text {out }}(t, x, \xi)\right) \\
& +\sum_{\ell=1}^{n} \int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s, \quad(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out, },}} \tag{77}
\end{align*}
$$

4) For $j \in\{m+1, \ldots, n\}$, integrating (74) along the characteristic curve $\chi_{i j}(s ; t, x, \xi)$ for $s \in$ $\left(t, s_{i j}^{\text {out }}(t, x, \xi)\right)$ yields the following integral equation:

$$
\begin{equation*}
k_{i j}(t, x, \xi)=\widetilde{k}_{i j}^{0}(t, x, \xi)+\sum_{\ell=1}^{n} \int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s \tag{78}
\end{equation*}
$$

5) We now want to plug $(78)$ into $(76)$ and $(77)$, respectively. From $(78)$ we have

$$
\begin{align*}
& k_{i, m+\ell}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)=\widetilde{k}_{i, m+\ell}^{0}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right) \\
&+\sum_{q=1}^{n} \int_{s_{i j}^{\text {out }}(t, x, \xi)}^{s_{i, m+\ell}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)} k_{i q}\left(\chi_{i, m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)\right) \\
& \times \widetilde{m}_{q, m+\ell}^{1}\left(s, \chi_{m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), 0\right)\right) d s \tag{79}
\end{align*}
$$

Plugging (79) into (76) and (77), we obtain, for every $j \in\{i, \ldots, m\}$ and $(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out, }}}$,

$$
\begin{aligned}
& k_{i j}(t, x, \xi)=\widehat{k}_{i j}^{0}(t, x, \xi) \\
& +\sum_{\ell=1}^{p}\left(\sum_{q=1}^{n} \int_{s_{i j}^{\text {out }}}^{s_{i, m+\ell}^{\text {out }}(t, x, \xi)} \begin{array}{l}
\text { out } \\
\left.i j, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right) \\
k_{i q}\left(\chi_{i, m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)\right) \\
\left.\times \widetilde{m}_{q, m+\ell}^{1}\left(s, \chi_{m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), 0\right)\right) d s\right) \widetilde{q}_{\ell j}^{1}\left(s_{i j}^{\text {out }}(t, x, \xi)\right) \\
\\
\quad+\sum_{\ell=1}^{n} \int_{t}^{s_{i j}^{\text {out }}(t, x, \xi)} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s,
\end{array}, l\right.
\end{aligned}
$$

where we introduced

$$
\widehat{k}_{i j}^{0}(t, x, \xi)=\sum_{\ell=1}^{p} \widetilde{k}_{i, m+\ell}^{0}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right) \widetilde{q}_{\ell j}^{1}\left(s_{i j}^{\text {out }}(t, x, \xi)\right)
$$

Note that

$$
\begin{equation*}
\widehat{k}_{i j}^{0} \in C^{0}\left(\overline{\mathcal{P}_{i j}^{\text {out,- }}}\right) \cap L^{\infty}\left(\mathcal{P}_{i j}^{\text {out },-}\right), \quad \forall j \in\{i+1, \ldots, m\} \tag{80}
\end{equation*}
$$

and, since $\mathcal{P}_{i i}^{\text {out, }-}=\mathcal{P}$ (because of (20p),

$$
\begin{equation*}
\widehat{k}_{i i}^{0} \in C^{0}(\overline{\mathcal{P}}) \cap L^{\infty}(\mathcal{P}) \tag{81}
\end{equation*}
$$

Remark 3.4. Observe that, in general, for $j \in\{i+1, \ldots, m\}$, we have

$$
\widehat{k}_{i j}^{0} \neq \widetilde{k}_{i j}^{0} \text { on } \mathcal{D}_{i j}^{\text {out }}
$$

### 3.2.4 Solution to the integral equations

564 In this subsection we show that there exists a unique solution to the system of integral equations of the previous section. This will conclude the proof of Theorem 2.6.

Fixed-point argument. As it is classical, we reformulate the existence of such a solution into the existence of a fixed-point of the mapping defined by the right-hand sides of these equations. Let us first introduce $K^{0}=\left(k_{i j}^{0}\right)_{\substack{1 \leq i \leq m \\ 1 \leq j \leq n}}$ defined by

$$
k_{i j}^{0}(t, x, \xi)= \begin{cases}\widetilde{k}_{i j}^{0}(t, x, \xi) & \text { if } j \in\{1, \ldots, i-1\} \\ \widehat{k}_{i i}^{0}(t, x, \xi) & \text { if } j=i, \\ \widetilde{k}_{i j}^{0}(t, x, \xi) & \text { if } j \in\{i+1, \ldots, m\} \text { and }(t, x, \xi) \in \mathcal{P}_{i j}^{\text {out },+}, \\ \widehat{k}_{i j}^{0}(t, x, \xi) & \text { if } j \in\{i+1, \ldots, m\} \text { and }(t, x, \xi) \in \mathcal{P}_{i j}^{\text {out },-} \\ \widetilde{k}_{i j}^{0}(t, x, \xi) & \text { if } j \in\{m+1, \ldots, n\}\end{cases}
$$

Thanks in particular to 85,80 and (81), we see that

$$
\begin{aligned}
& k_{i j}^{0} \in C^{0}\left(\overline{\mathcal{P}_{i j}^{\text {out,+ }}}\right) \cap C^{0}\left(\overline{\mathcal{P}_{i j}^{\text {out },-}}\right) \cap L^{\infty}(\mathcal{P}) \quad \text { if } j \in\{i+1, \ldots, m\}, \\
& k_{i j}^{0} \in C^{0}(\overline{\mathcal{P}}) \cap L^{\infty}(\mathcal{P}) \quad \text { otherwise. }
\end{aligned}
$$

It is this regularity that dictates the space in which we can work. More precisely, let us introduce the vector space $B$ defined by

$$
B=\left\{\begin{array}{ll} 
& k_{i j} \in C^{0}\left(\overline{\mathcal{P}_{i j}^{\text {out,+ }}}\right) \cap C^{0}\left(\overline{\mathcal{P}_{i j}^{\text {out, }-}}\right) \cap L^{\infty}(\mathcal{P}) \text { if } j \in\{i+1, \ldots, m\},  \tag{82}\\
K=\left(k_{i j}\right)_{\substack{1 \leq i \leq m \\
1 \leq j \leq n}} & k_{i j} \in C^{0}(\overline{\mathcal{P}}) \cap L^{\infty}(\mathcal{P}) \text { otherwise. }
\end{array}\right\} .
$$

We can check that $B$ is a Banach space when equipped with the $L^{\infty}$ norm. Let us now introduce the mapping

$$
\Phi: B \longrightarrow B
$$

defined, for every $K \in B$, by

$$
\Phi(K)=K^{0}+\Phi_{1}(K)+\Phi_{2}(K)
$$

where, for every $(t, x, \xi) \in \overline{\mathcal{P}}$,

$$
\begin{align*}
& \left(\Phi_{1}(K)\right)_{i j}(t, x, \xi)= \\
& \begin{cases}-\sum_{\ell=1}^{n} \int_{s_{s_{i j}^{i j}}^{t}(t, x, \xi)}^{t} k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s, & \text { if } j \in\{1, \ldots, i-1\}, \\
\sum_{\ell=1}^{n} \int_{t}^{s_{i j}^{\text {iit }}(t, x, \xi)} & k_{i \ell}\left(\chi_{i j}(s ; t, x, \xi)\right) \widetilde{m}_{\ell j}^{1}\left(s, \chi_{j}(s ; t, \xi)\right) d s, \\
\text { if } j \in\{i, \ldots, n\},\end{cases} \tag{83}
\end{align*}
$$

and

$$
\begin{align*}
& \left(\Phi_{2}(K)\right)_{i j}(t, x, \xi)= \\
& \sum_{\ell=1}^{p}\left(\sum_{q=1}^{n} \int_{s_{i j}^{\text {out }}(t, x, \xi)}^{s_{i, m}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)} k_{i q}\left(\chi_{i, m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)\right)\right. \\
& \left.\times \widetilde{m}_{q, m+\ell}^{1}\left(s, \chi_{m+\ell}\left(s ; s_{i j}^{\text {out }}(t, x, \xi), 0\right)\right) d s\right) \widetilde{q}_{\ell j}^{1}\left(s_{i j}^{\text {out }}(t, x, \xi)\right), \tag{84}
\end{align*}
$$

Contraction of the mapping. We will now prove that $\Phi^{N}$ is a contraction for $N \in \mathbb{N}^{*}$ large enough. Therefore, the Banach fixed-point theorem can be applied, giving the existence (and uniqueness) of $K \in B$ such that

$$
K=\Phi(K) .
$$

This will conclude the proof of Theorem 2.6. Now, to show that $\Phi^{N}$ is a contraction when $N$ is large, it is sufficient to prove the following estimate:

Proposition 3.6. There exists $C>0$ such that, for every $N \in \mathbb{N}^{*}$ and $K, H \in B$,

$$
\begin{equation*}
\left\|\Phi^{N}(K)-\Phi^{N}(H)\right\|_{L^{\infty}(\mathcal{P})^{m \times n}} \leq \frac{C^{N}}{N!}\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}} \tag{86}
\end{equation*}
$$

To establish 86 we will use the following key lemma:
Lemma 3.7. For every $i \in\{1, \ldots, m\}$, there exist a function $\Omega_{i} \in C^{1}(\overline{\mathcal{P}}) \cap L^{\infty}(\mathcal{P})$ and $\varepsilon_{0}>0$ such that, for every $(t, x, \xi) \in \overline{\mathcal{P}}$, we have $\Omega_{i}(t, x, \xi) \geq 0$ with

$$
\begin{equation*}
\frac{\partial \Omega_{i}}{\partial t}(t, x, \xi)+\lambda_{i}(t, x) \frac{\partial \Omega_{i}}{\partial x}(t, x, \xi)+\lambda_{j}(t, \xi) \frac{\partial \Omega_{i}}{\partial \xi}(t, x, \xi) \geq \varepsilon_{0}, \quad \forall j \in\{1, \ldots, i-1\} \tag{87}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial \Omega_{i}}{\partial t}(t, x, \xi)+\lambda_{i}(t, x) \frac{\partial \Omega_{i}}{\partial x}(t, x, \xi)+\lambda_{j}(t, \xi) \frac{\partial \Omega_{i}}{\partial \xi}(t, x, \xi) \leq-\varepsilon_{0}, \quad \forall j \in\{i, \ldots, n\} \tag{88}
\end{equation*}
$$

The proof of Lemma 3.7 is postponed to Appendix Cor the sake of the presentation.
Remark 3.8. In the time-independent case, we can take $\Omega_{i}(x, \xi)=\phi_{i}(x)-\nu \phi_{i}(\xi)$ (we recall that $\left.\phi_{i}(x)=\int_{0}^{x} \frac{1}{-\lambda_{i}(y)} d y\right)$ where $\nu \in[0,1)$ is any number such that $\nu>\max _{1 \leq j<i \leq m} \max _{\xi \in[0,1]} \lambda_{i}(\xi) / \lambda_{j}(\xi)$. This function appeared for instance in HDMVK16, Lemma 6.2] for systems with constant coefficients and in HVDMK19, (A.32)] for systems with time-independent coefficients (see also [CVKB13, Lemma A.4] for $2 \times 2$ systems, where it is enough to take $\nu=0$ since 87) becomes void).

Remark 3.9. Observe that it follows from the estimate (88) that, for every $j \in\{i, \ldots, n\}$,

$$
s \mapsto \Omega_{i}\left(\chi_{i j}(s ; t, x, \xi)\right) \text { is strictly decreasing. }
$$

598 This is the analogue to HDMVK16, Remark 10].
We can now prove Proposition 3.6
Proof of Proposition 3.6. Let us denote by

$$
R=\max \left\{\left\|\widetilde{M}^{1}\right\|_{L^{\infty}(\mathbb{R} \times(0,1))^{n \times n}},\left\|\widetilde{Q}^{1}\right\|_{L^{\infty}(\mathbb{R})^{p \times m}}\right\}
$$

1) We start with the estimate of $\left\|\Phi_{1}(K)-\Phi_{1}(H)\right\|_{L^{\infty}(\mathcal{P})^{m \times n}}$. Set

$$
C_{1}=\frac{n}{\varepsilon_{0}} R .
$$

Let $j \in\{1, \ldots, i-1\}$. From the definition 83 of $\Phi_{1}$ we see that

$$
\left|\left(\Phi_{1}(K)-\Phi_{1}(H)\right)_{i j}(t, x, \xi)\right| \leq n R\left(\int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} 1 d s\right)\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
$$

Thanks to the estimate (87) we can perform the change of variable $s \mapsto \theta(s)=\Omega_{i}\left(\chi_{i j}(s ; t, x, \xi)\right)$ and obtain

$$
\varepsilon_{0}\left(\int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} 1 d s\right) \leq \int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} \frac{d \theta}{d s}(s) d s=\theta(t)-\theta\left(s_{i j}^{\mathrm{in}}(t, x, \xi)\right) \leq \theta(t)=\Omega_{i}(t, x, \xi)
$$

This gives the estimate

$$
\left|\left(\Phi_{1}(K)-\Phi_{1}(H)\right)_{i j}(t, x, \xi)\right| \leq C_{1} \Omega_{i}(t, x, \xi)\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
$$

It is important to point out that the right-hand side does not depend on the second index $j$. Computing $\Phi_{1}^{2}(H)-\Phi_{1}^{2}(K)=\Phi_{1}\left(\Phi_{1}(H)\right)-\Phi_{1}\left(\Phi_{1}(K)\right)$ and using the previous estimate, we obtain

$$
\begin{aligned}
& \left|\left(\Phi_{1}^{2}(K)-\Phi_{1}^{2}(H)\right)_{i j}(t, x, \xi)\right| \\
& \quad \leq n R C_{1}\left(\int_{s_{i j}^{\mathrm{in}(t, x, \xi)}}^{t} \Omega_{i}\left(\chi_{i j}(s ; t, x, \xi)\right) d s\right)\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
\end{aligned}
$$

Using again the change of variable $s \mapsto \theta(s)$ and (87), we obtain

$$
\begin{aligned}
& \varepsilon_{0}\left(\int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} \Omega_{i}\left(\chi_{i j}(s ; t, x, \xi)\right) d s\right)=\varepsilon_{0} \int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} \theta(s) d s \leq \int_{s_{i j}^{\mathrm{in}}(t, x, \xi)}^{t} \theta(s) \frac{d \theta}{d s}(s) d s \\
&=\frac{\theta(t)^{2}}{2}-\frac{\theta\left(s_{i j}^{\mathrm{in}}(t, x, \xi)\right)^{2}}{2} \leq \frac{\theta(t)^{2}}{2}=\frac{\Omega_{i}(t, x, \xi)^{2}}{2}
\end{aligned}
$$

This gives the estimate

$$
\left|\left(\Phi_{1}^{2}(K)-\Phi_{1}^{2}(H)\right)_{i j}(t, x, \xi)\right| \leq \frac{C_{1}^{2} \Omega_{i}(t, x, \xi)^{2}}{2}\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
$$

By induction, we easily obtain that, for every $N \in \mathbb{N}^{*}$,

$$
\begin{equation*}
\left|\left(\Phi_{1}^{N}(K)-\Phi_{1}^{N}(H)\right)_{i j}(t, x, \xi)\right| \leq \frac{C_{1}^{N} \Omega_{i}(t, x, \xi)^{N}}{N!}\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}} . \tag{89}
\end{equation*}
$$

Using the estimate (88) instead of (87), we can obtain exactly the same estimate as 89) for $j \in\{i, \ldots, n\}$. Since $\Omega_{i}$ is bounded, it follows that

$$
\left\|\Phi_{1}^{N}(K)-\Phi_{1}^{N}(H)\right\|_{L^{\infty}(\mathcal{P})^{m \times n}} \leq \frac{C^{N}}{N!}\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}},
$$

for some $C$ independent of $N$ and $K, H$.
2) Let us now take care of $\Phi_{2}(K)-\Phi_{2}(H)$. The idea to estimate this term is essentially the same as before, with the extra use of the decreasing property stated in Remark 3.9. Set

$$
C_{2}=\frac{n}{\varepsilon_{0}} R^{2} p
$$

From the definition (84) of $\Phi_{2}$ we see that, for $j \in\{i, \ldots, m\}$ and $(t, x, \xi) \in \overline{\mathcal{P}_{i j}^{\text {out,-- }}}$,

$$
\begin{aligned}
& \left|\left(\Phi_{2}(K)-\Phi_{2}(H)\right)_{i j}(t, x, \xi)\right| \\
& \quad \leq n R^{2} \sum_{\ell=1}^{p}\left(\int_{s_{i j}^{\text {out }}(t, x, \xi)}^{s_{i, m}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)} 1 d s\right)\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}} .
\end{aligned}
$$

Thanks to the estimate (88) we can perform again the change of variable

$$
s \mapsto \theta(s)=\Omega_{i}\left(\chi_{i j}(s ; t, x, \xi)\right),
$$

which is decreasing since $j \geq i$ (see Remark 3.9), and obtain

$$
\begin{aligned}
& \varepsilon_{0}\left(\int_{s_{i j}^{\text {out }}(t, x, \xi)}^{s_{i, \xi}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)} 1 d s\right) \\
& \quad \leq \int_{s_{i j}^{\text {out }}(t, x, \xi)}^{s_{i, m}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, x) ; t, x\right), 0\right)}-\frac{d \theta}{d s}(s) d s \\
& \quad=-\theta\left(s_{i, m+\ell}^{\text {out }}\left(s_{i j}^{\text {out }}(t, x, \xi), \chi_{i}\left(s_{i j}^{\text {out }}(t, x, \xi) ; t, x\right), 0\right)\right)+\theta\left(s_{i j}^{\text {out }}(t, x, \xi)\right) \leq \theta(t)=\Omega_{i}(t, x, \xi)
\end{aligned}
$$

This gives the estimate

$$
\left|\left(\Phi_{2}(K)-\Phi_{2}(H)\right)_{i j}(t, x, \xi)\right| \leq C_{2} \Omega_{i}(t, x, \xi)\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
$$

Note that this estimate is also valid if $j \notin\{i, \ldots, m\}$ or $(t, x, \xi) \notin \overline{\mathcal{P}_{i j}^{\text {out,-- }}}$ since $\left(\Phi_{2}(\cdot)\right)_{i j}=0$ in this case. Reasoning by induction as before, it is now not difficult to obtain the estimate

$$
\left\|\Phi_{2}^{N}(K)-\Phi_{2}^{N}(H)\right\|_{L^{\infty}(\mathcal{P})^{m \times n}} \leq \frac{C^{N}}{N!}\|K-H\|_{L^{\infty}(\mathcal{P})^{m \times n}}
$$

for some $C$ independent of $N$ and $K, H$.

### 3.3 On the time-periodicity of $F$

In this section we assume that, for some $\tau>0, \Lambda, M$ and $Q$ are $\tau$-periodic with respect to time and show that the above construction of $F$ leads, with minor modifications, to a $F$ which is also $\tau$-periodic with respect to time.

First of all, concerning the extension of $\Lambda$ to a function of $\mathbb{R}^{2}$ (and of $M$ and $Q$ whenever needed), it is clear that one can extend $\Lambda$ to $\mathbb{R} \times[0,1]$ by just requiring the $\tau$-periodicity with respect to time of this extension. This extension is still denoted by $\Lambda$. Then one extends $\Lambda$ to $\mathbb{R}^{2}$ so that this extension, still denoted by $\Lambda$, is $\tau$-periodic with respect to time and so that the properties (2), (3), (4) and (5) remain valid on $\mathbb{R}^{2}$ (see e.g. Remark 1.4 ).

From the construction of $F$ (see (30), (38) and (50)) it is clear that $F$ is $\tau$-periodic with respect to time if so are all the matrix-valued functions involved in the several transformations of this article. Now, in order to obtain the $\tau$-periodicity of these transformations, the minor modifications/comments are essentially the following ones.

1) Concerning the diagonal transformation to remove the diagonal terms in $M$ (see Section 2.1), one simply observes that the function (33) is $\tau$-periodic with respect to time if so is $m_{i i}$, thanks to the properties

$$
\begin{equation*}
\chi_{i}(s+\tau ; t+\tau, \xi)=\chi_{i}(s ; t, \xi), \quad s_{i}^{\mathrm{in}}(t+\tau, \xi)=s_{i}^{\mathrm{in}}(t, \xi)+\tau \tag{90}
\end{equation*}
$$

2) Concerning the kernel $H$ of the Fredholm transformation (see Section 3.1) one easily checks that it is indeed $\tau$-periodic with respect to time. This follows from the uniqueness of the solution to (59) and similar properties to (90).
3) Concerning the kernel $K$ of the Volterra transformation of the second kind (see Section 3.2), to construct it in such a way that it is $\tau$-periodic with respect to time, it suffices to observe that $\widetilde{m}_{\ell j}^{1}, r_{i j}$ and $q_{\ell j}^{1}$ become $\tau$-periodic with respect to time once $M$ and $Q$ are, and to modify the definition of the space $B$ given in (82) by adding the condition that the $k_{i j}, 1 \leq i \leq m$, $1 \leq j \leq n$, are $\tau$-periodic with respect to time (alternatively, one can keep (82) and deduce from the uniqueness of the fixed point of $\Phi$ that it has to be $\tau$-periodic with respect to time).

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## A Background on broad solutions

We recall that all the systems of this paper have the following form:

$$
\left\{\begin{array}{l}
\frac{\partial y}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial y}{\partial x}(t, x)=M(t, x) y(t, x)+G(t, x) y(t, 0)  \tag{91}\\
y_{-}(t, 1)=\int_{0}^{1} F(t, \xi) y(t, \xi) d \xi, \quad y_{+}(t, 0)=Q(t) y_{-}(t, 0) \\
y\left(t^{0}, x\right)=y^{0}(x)
\end{array}\right.
$$

where $M$ and $Q$ have at least the regularity (5), $F \in L^{\infty}((0,+\infty) \times(0,1))^{m \times n}$ and

$$
G \in C^{0}([0,+\infty) \times[0,1])^{n \times n} \cap L^{\infty}((0,+\infty) \times(0,1))^{n \times n}
$$

## 641 A. 1 Definition of broad solution

Let us now introduce the notion of solution for such systems. To this end, we have to restrict our discussion to the domain where the system (91) evolves, i.e. on $\left(t^{0},+\infty\right) \times(0,1)$. For every $(t, x) \in\left(t^{0},+\infty\right) \times(0,1)$, we have

$$
\left(s, \chi_{i}(s ; t, x)\right) \in\left(t^{0},+\infty\right) \times(0,1), \quad \forall s \in\left(\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right), s_{i}^{\text {out }}(t, x)\right)
$$

where we introduced

$$
\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)=\max \left\{t^{0}, s_{i}^{\mathrm{in}}(t, x)\right\}<t
$$

Formally, writing the $i$-th equation of the system (91) along the characteristic $\chi_{i}(s ; t, x)$ for $s \in$ $\left[\bar{s}_{i}^{\text {in }}\left(t^{0} ; t, x\right), s_{i}^{\text {out }}(t, x)\right]$, and using the chain rules yields the ODE

$$
\left\{\begin{array}{l}
\frac{d}{d s} y_{i}\left(s, \chi_{i}(s ; t, x)\right)=\sum_{j=1}^{n} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right)+\sum_{j=1}^{n} g_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}(s, 0),  \tag{92}\\
y_{i}\left(\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right), \chi_{i}\left(\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right) ; t, x\right)\right)=b_{i}(y)(t, x)
\end{array}\right.
$$

where the initial condition $b_{i}(y)(t, x)$ is given by the appropriate boundary or initial conditions of the system (91):

$$
b_{i}(y)(t, x)= \begin{cases}\sum_{j=1}^{n} \int_{0}^{1} f_{i j}\left(s_{i}^{\text {in }}(t, x), \xi\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) d \xi & \text { if } s_{i}^{\text {in }}(t, x)>t^{0} \text { and } i \in\{1, \ldots, m\}  \tag{93}\\ \sum_{j=1}^{m} q_{i-m, j}\left(s_{i}^{\text {in }}(t, x)\right) y_{j}\left(s_{i}^{\text {in }}(t, x), 0\right) & \text { if } s_{i}^{\text {in }}(t, x)>t^{0} \text { and } i \in\{m+1, \ldots, n\}, \\ y_{i}^{0}\left(\chi_{i}\left(t^{0} ; t, x\right)\right) & \text { if } s_{i}^{\text {in }}(t, x)<t^{0}\end{cases}
$$

Integrating the ODE (92) over $s \in\left[\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right), t\right]$, we obtain the following system of integral equations:

$$
\begin{align*}
& y_{i}(t, x)=b_{i}(y)(t, x)+\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s \\
&+\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}\left(t^{0} ; t, x\right)}}^{t} g_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}(s, 0) d s \tag{94}
\end{align*}
$$

Definition A.1. Let $t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$ be fixed. We say that a function $y:\left(t^{0},+\infty\right) \times$

$$
\begin{equation*}
y \in C^{0}\left(\left[t^{0}, t^{0}+T\right] ; L^{2}(0,1)^{n}\right) \cap C^{0}\left([0,1] ; L^{2}\left(t^{0}, t^{0}+T\right)^{n}\right), \quad \forall T>0 \tag{95}
\end{equation*}
$$

and if the integral equation (94) is satisfied for every $i \in\{1, \ldots, n\}$, for a.e. $t>t^{0}$ and a.e. $x \in(0,1)$.

## A. 2 Well-posedness

This section is devoted to the following well-posedness result regarding system (91):
Theorem A.2. For every $t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$, there exists a unique broad solution to (91). Moreover, there exists $C>0$ such that, for every $T>0, t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$, the corresponding broad solution y satisfies

$$
\begin{equation*}
\|y\|_{C^{0}\left(\left[t^{0}, t^{0}+T\right] ; L^{2}(0,1)^{n}\right)}+\|y\|_{C^{0}\left([0,1] ; L^{2}\left(t^{0}, t^{0}+T\right)^{n}\right)} \leq C e^{C T}\left\|y^{0}\right\|_{L^{2}(0,1)^{n}} \tag{96}
\end{equation*}
$$

Remark A.3. It follows from the uniformity of the constant $C e^{C T}$ with respect to the initial time $t^{0}$ in the estimate (96) that, for systems of the form (91), the uniform stability property (8) is a consequence of the finite-time global attractor property (7) (simply take $\delta>0$ such that $C e^{C T} \delta \leq \varepsilon$ ).

Let us first point out that this well-posedness result for our initial system (1) for the particular $F$ that we have constructed in Section 2 follows in fact from the well-posedness result for the final target system of Proposition 2.11 (easier to establish), since we have shown that both systems are equivalent by means of several invertible transformations. However, it is still important to have such a well-posedness result for any $F$ within the class studied, which is a result that also has its own interest. We will provide a complete proof since, to the best of our knowledge, there are no references that show the well-posedness for the initial-boundary value problem (91) with non-local terms $G(t, x) y(t, 0)$, with weak regularity (95) and with uniform estimate (96).

Proof of Theorem A.2. We first remark that it is enough to prove the theorem for $\|Q\|_{L^{\infty}}$ small enough, say

$$
\begin{equation*}
\|Q\|_{L^{\infty}} \leq \alpha \tag{97}
\end{equation*}
$$

where $\alpha>0$ does not depend on $T, t^{0}, y^{0}$ nor on $M, G, F$. This follows from the following change of variable:

$$
y=D \tilde{y}, \quad D=\left(\begin{array}{cc}
\frac{\alpha}{\|Q\|_{L \infty}+\alpha} \mathrm{Id}_{\mathbb{R}^{m}} & 0 \\
0 & \mathrm{Id}_{\mathbb{R}^{p}}
\end{array}\right)
$$

where $\tilde{y}$ is the solution to the system $(\tilde{M}, \tilde{G}, \tilde{F}, \tilde{Q})$ with

$$
\tilde{M}=D^{-1} M D, \quad \tilde{G}=D^{-1} G D, \quad \tilde{F}=\left(\frac{\alpha}{\|Q\|_{L^{\infty}}+\alpha}\right)^{-1} F D, \quad \tilde{Q}=\frac{\alpha}{\|Q\|_{L^{\infty}}+\alpha} Q
$$

Let us now show how to prove the theorem under the smallness condition 97 with the Banach fixed point theorem ( $\alpha>0$ will be fixed adequately below). Let $T>0, t^{0} \geq 0$ and $y^{0} \in L^{2}(0,1)^{n}$ be fixed for the remainder of the proof. It is clear that a function $y:\left(t^{0},+\infty\right) \times(0,1) \longrightarrow \mathbb{R}^{n}$ satisfies the integral equation 94 if, and only if, it is a fixed point of the map $\mathcal{F}: B \longrightarrow B$, where

$$
B=C^{0}\left(\left[t^{0}, t^{0}+T\right] ; L^{2}(0,1)^{n}\right) \cap C^{0}\left([0,1] ; L^{2}\left(t^{0}, t^{0}+T\right)^{n}\right),
$$

and $(\mathcal{F}(y))_{i}(t, x)$ is given by the expression on the right-hand side of (94). It can be checked that $\mathcal{F}$ indeed maps $B$ into itself (actually, by computations similar to the upcoming ones). Let us now make $B$ a Banach space by equipping it with the following weighted norm:

$$
\|y\|_{B}=\|y\|_{B_{1}}+\|y\|_{B_{2}},
$$

where

$$
\|y\|_{B_{1}}=\max _{t \in\left[t^{0}, t^{0}+T\right]} e^{-\frac{L_{1}}{2}\left(t-t^{0}\right)} \sqrt{\int_{0}^{1} \sum_{i=1}^{n}\left|y_{i}(t, x)\right|^{2} e^{-L_{2} x} d x}
$$

and

$$
\|y\|_{B_{2}}=\max _{x \in[0,1]} e^{\frac{L_{2}}{2}(1-x)} \sqrt{\int_{t^{0}}^{t^{0}+T} \sum_{i=1}^{n}\left|y_{i}(t, x)\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t}
$$

where $L_{1}, L_{2}>0$ are constants independent of $T, t^{0}$ and $y^{0}$ that will be fixed below. Our goal is to show that, for $L_{1}, L_{2}>0$ large enough,

$$
\begin{equation*}
\left\|\mathcal{F}\left(y^{1}\right)-\mathcal{F}\left(y^{2}\right)\right\|_{B} \leq \frac{1}{2}\left\|y^{1}-y^{2}\right\|_{B}, \quad \forall y^{1}, y^{2} \in B \tag{98}
\end{equation*}
$$

It is then not difficult to check that the fixed point of $\mathcal{F}$ satisfies the estimate (96). Indeed, using (98), we easily see that the fixed point $y$ of $\mathcal{F}$ will satisfy

$$
\frac{1}{2}\|y\|_{B} \leq\|\mathcal{F}(0)\|_{B}
$$

and some straightforward computations show that

$$
\begin{gathered}
\|y\|_{C^{0}\left(\left[t^{0}, t^{0}+T\right] ; L^{2}(0,1)^{n}\right)}^{2} \leq e^{L_{2}} e^{L_{1} T}\|y\|_{B_{1}}^{2}, \quad\|y\|_{C^{0}\left([0,1] ; L^{2}\left(t^{0}, t^{0}+T\right)^{n}\right)}^{2} \leq e^{L_{1} T}\|y\|_{B_{2}}^{2}, \\
\|\mathcal{F}(0)\|_{B}^{2} \leq 2\left(1+\frac{e^{L_{2}}}{\varepsilon}\right) e^{\left\|\frac{\partial \Lambda}{\partial x}\right\|_{L^{\infty}}^{T}}\left\|y^{0}\right\|_{L^{2}(0,1)^{n}}^{2} .
\end{gathered}
$$

Let us now establish (98). We introduce

$$
y=y^{1}-y^{2}
$$

so that $\mathcal{F}\left(y^{1}\right)-\mathcal{F}\left(y^{2}\right)$ is equal to the right-hand side of 94 with $y^{0}=0$. We have to estimate four types of terms in each $\|\cdot\|_{B_{i}}$-norm $(i=1,2)$. For convenience, we denote by

$$
R_{1}=\max \left\{\|\Lambda\|_{L^{\infty}},\left\|\frac{\partial \Lambda}{\partial x}\right\|_{L^{\infty}}\right\}, \quad R_{2}=\max \left\{\|M\|_{L^{\infty}},\|G\|_{L^{\infty}},\|F\|_{L^{\infty}}\right\}
$$

We recall that it is crucial that $\alpha$ does not depend on $R_{2}$.
Estimate of the $\|\cdot\|_{B_{1}}$-norm. Let $t \in\left[t^{0}, t^{0}+T\right]$ be fixed. Let $I=\left\{x \in(0,1), \quad s_{i}^{\text {in }}(t, x)>t^{0}\right\}$.

1) Let $i \in\{1, \ldots, m\}$. For a.e. $x \in I$, using Cauchy-Schwarz inequality, we have

$$
\left|\sum_{j=1}^{n} \int_{0}^{1} f_{i j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) d \xi\right|^{2} \leq n R_{2}^{2} e^{L_{2}}\|y\|_{B_{1}}^{2} e^{L_{1}\left(s_{i}^{\mathrm{in}}(t, x)-t^{0}\right)}
$$

Using a finer version of (15), namely,

$$
\begin{equation*}
\frac{1-x}{R_{1}} \leq t-s_{i}^{\mathrm{in}}(t, x) \tag{99}
\end{equation*}
$$

(obtained similarly to 19 ) we obtain the estimate

$$
\int_{I}\left|\sum_{j=1}^{n} \int_{0}^{1} f_{i j}\left(s_{i}^{\text {in }}(t, x), \xi\right) y_{j}\left(s_{i}^{\text {in }}(t, x), \xi\right) d \xi\right|^{2} e^{-L_{2} x} d x \leq\left(n R_{2}^{2} \frac{1}{\frac{L_{1}}{R_{1}}-L_{2}}\right) e^{L_{1}\left(t-t^{0}\right)}\|y\|_{B_{1}}^{2},
$$

provided that

$$
\begin{equation*}
\frac{L_{1}}{R_{1}}-L_{2}>0 \tag{100}
\end{equation*}
$$

2) Let $i \in\{m+1, \ldots, n\}$. Using (15), we have

$$
\begin{aligned}
\int_{I}\left|\sum_{j=1}^{m} q_{i-m, j}\left(s_{i}^{\mathrm{in}}(t, x)\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), 0\right)\right|^{2} & e^{-L_{2} x} d x \\
& \leq m \alpha^{2} e^{L_{1}\left(t-t^{0}\right)} \sum_{j=1}^{m} \int_{I}\left|y_{j}\left(s_{i}^{\mathrm{in}}(t, x), 0\right)\right|^{2} e^{-L_{1}\left(s_{i}^{\mathrm{in}}(t, x)-t^{0}\right)} d x
\end{aligned}
$$

Doing the change of variables $\sigma=s_{i}^{\text {in }}(t, x)$ and using the estimate (see (14), 12) and (19) )

$$
\frac{\partial s_{i}^{\operatorname{in}}(t, x)}{\partial x}=\frac{-e^{-\int_{s_{i}^{\operatorname{in}(t, x)}}^{t} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta}}{\lambda_{i}\left(s_{i}^{\operatorname{in}}(t, x), 0\right)} \leq-\frac{1}{R_{1}} e^{-\frac{R_{1}}{\varepsilon}}
$$

we obtain

$$
\int_{I}\left|\sum_{j=1}^{m} q_{i-m, j}\left(s_{i}^{\mathrm{in}}(t, x)\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), 0\right)\right|^{2} e^{-L_{2} x} d x \leq\left(m \alpha^{2} R_{1} e^{\frac{R_{1}}{\varepsilon}} e^{-L_{2}}\right) e^{L_{1}\left(t-t^{0}\right)}\|y\|_{B_{2}}^{2}
$$

3) For the next term, we have

$$
\begin{aligned}
& \int_{0}^{1} \sum_{i=1}^{n}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\text {in }}\left(t^{0} ; t, x\right)}^{t} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s\right|^{2} e^{-L_{2} x} d x \\
& \leq n R_{2}^{2} \frac{1}{\varepsilon} \sum_{i=1}^{n} \int_{0}^{1} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} \sum_{j=1}^{n}\left|y_{j}\left(s, \chi_{i}(s ; t, x)\right)\right|^{2} d s d x
\end{aligned}
$$

Using the change of variable $(\sigma, \xi)=\left(s, \chi_{i}(s ; t, x)\right)$, whose Jacobian determinant is (see 12 )

$$
\operatorname{det}\left(\begin{array}{cc}
1 & 0 \\
\lambda_{i}\left(s, \chi_{i}(s ; t, x)\right) & \frac{\partial \chi_{i}}{\partial x}(s ; t, x)
\end{array}\right)=e^{-\int_{s}^{t} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta} \geq e^{-\frac{R_{1}}{\varepsilon}}, \quad \forall s \in\left(s_{i}^{\text {in }}(t, x), t\right)
$$

we obtain

$$
\begin{aligned}
& \int_{0}^{1} \sum_{i=1}^{n}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s\right|^{2} e^{-L_{2} x} d x \\
& \leq\left(n^{2} R_{2}^{2} \frac{1}{\varepsilon} e^{\frac{R_{1}}{\varepsilon}} e^{L_{2}} \frac{1}{L_{1}}\right) e^{L_{1}\left(t-t^{0}\right)}\|y\|_{B_{1}}^{2}
\end{aligned}
$$

4) Finally, the estimate of the remaining term is easy:

$$
\begin{array}{rl}
\int_{0}^{1} \sum_{i=1}^{n}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} g_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}(s, 0) d s\right|^{2} e^{-L_{2} x} & d x \\
& \leq\left(n^{2} R_{2}^{2} \frac{1}{\varepsilon} e^{-L_{2}}\right) e^{L_{1}\left(t-t^{0}\right)}\|y\|_{B_{2}}^{2}
\end{array}
$$

In summary, we have established the following estimate (provided that (100) holds):

$$
\begin{align*}
\left\|\mathcal{F}\left(y^{1}\right)-\mathcal{F}\left(y^{2}\right)\right\|_{B_{1}}^{2} \leq 3\left(m n R_{2}^{2} \frac{1}{\frac{L_{1}}{R_{1}}-L_{2}}\right. & \left.+n^{2} R_{2}^{2} \frac{1}{\varepsilon} e^{\frac{R_{1}}{\varepsilon}} e^{L_{2}} \frac{1}{L_{1}}\right)\|y\|_{B_{1}}^{2} \\
& +3\left((n-m) m \alpha^{2} R_{1} e^{\frac{R_{1}}{\varepsilon}} e^{-L_{2}}+n^{2} R_{2}^{2} \frac{1}{\varepsilon} e^{-L_{2}}\right)\|y\|_{B_{2}}^{2} \tag{101}
\end{align*}
$$

Estimate of the $\|\cdot\|_{B_{2}}$-norm. Let $x \in[0,1]$ be fixed. Let $J=\left\{t \in\left(t^{0}, t^{0}+T\right), \quad s_{i}^{\text {in }}(t, x)>t^{0}\right\}$.

1) Let $i \in\{1, \ldots, m\}$. We have

$$
\begin{aligned}
& \int_{J}\left|\sum_{j=1}^{n} \int_{0}^{1} f_{i j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) d \xi\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \\
& \leq n R_{2}^{2} \int_{0}^{1}\left(\int_{J} \sum_{j=1}^{n}\left|y_{j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right)\right|^{2} e^{-L_{1}\left(s_{i}^{\mathrm{in}}(t, x)-t^{0}\right)} e^{L_{1}\left(s_{i}^{\mathrm{in}}(t, x)-t\right)} d t\right) d \xi
\end{aligned}
$$

Using once again (99), 100), performing the change of variable $\sigma=s_{i}^{\text {in }}(t, x)$, and using the estimate (see 14, 12) and (19))

$$
\begin{equation*}
\frac{\partial s_{i}^{\operatorname{in}}(t, x)}{\partial t}=\frac{\lambda_{i}(t, x) e^{-\int_{s_{i}^{\operatorname{in}(t, x)}}^{t} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta}}{\lambda_{i}\left(s_{i}^{\operatorname{in}}(t, x), 1\right)} \geq \frac{\varepsilon}{R_{1}} e^{-\frac{R_{1}}{\varepsilon}} \tag{102}
\end{equation*}
$$

we obtain the estimate

$$
\int_{J}\left|\sum_{j=1}^{n} \int_{0}^{1} f_{i j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) y_{j}\left(s_{i}^{\mathrm{in}}(t, x), \xi\right) d \xi\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \leq\left(n \frac{R_{1}}{\varepsilon} R_{2}^{2} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{2}}\right) e^{-L_{2}(1-x)}\|y\|_{B_{2}}^{2}
$$

2) The next estimate is where we will need the smallness assumption on $Q$. Let $i \in\{m+1, \ldots, n\}$. Using (15) and the change of variables $\sigma=s_{i}^{\text {in }}(t, x)$ (recall the estimate 102), we obtain

$$
\int_{J}\left|\sum_{j=1}^{m} q_{i-m, j}\left(s_{i}^{\text {in }}(t, x)\right) y_{j}\left(s_{i}^{\text {in }}(t, x), 0\right)\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \leq\left(m \alpha^{2} \frac{R_{1}}{\varepsilon} e^{\frac{R_{1}}{\varepsilon}}\right) e^{-L_{2}(1-x)}\|y\|_{B_{1}}^{2}
$$

3) For $i \in\{1, \ldots, m\}$, using the estimate

$$
-R_{1}(t-s) \leq x-\chi_{i}(s ; t, x)
$$

we have

$$
\begin{aligned}
\mid \sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} m_{i j}(s, & \left.\chi_{i}(s ; t, x)\right)\left.y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} \\
& \leq n R_{2}^{2} \frac{1}{\varepsilon} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} \sum_{j=1}^{n}\left|y_{j}\left(s, \chi_{i}(s ; t, x)\right)\right|^{2} e^{-L_{1}\left(s-t^{0}\right)} e^{-\frac{L_{1}}{R_{1}\left(\chi_{i}(s ; t, x)-x\right)}} d s
\end{aligned}
$$

Integrating and using the change of variable $(\sigma, \xi)=\left(s, \chi_{i}(s ; t, x)\right)$, whose Jacobian determinant is uniformly estimated for $s \in\left(s_{i}^{\text {in }}(t, x), t\right)$ by (see 12) and 19$)$

$$
\left|\operatorname{det}\left(\begin{array}{cc}
1 & 0 \\
\lambda_{i}\left(s, \chi_{i}(s ; t, x)\right) & \frac{\partial \chi_{i}}{\partial t}(s ; t, x)
\end{array}\right)\right|=\left|\lambda_{i}(t, x)\right| e^{-\int_{s}^{t} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}(\theta ; t, x)\right) d \theta} \geq \varepsilon e^{-\frac{R_{1}}{\varepsilon}},
$$

we obtain (using also that $x \leq \chi_{i}(s ; t, x) \leq 1$ )

$$
\begin{array}{r}
\int_{t^{0}}^{t^{0}+T} \sum_{i=1}^{m}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \\
\leq m n R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \int_{x}^{1}\left(\int_{t^{0}}^{t^{0}+T} \sum_{j=1}^{n}\left|y_{j}(\sigma, \xi)\right|^{2} e^{-L_{1}\left(\sigma-t^{0}\right)} d \sigma\right) e^{-\frac{L_{1}}{R_{1}}(\xi-x)} d \xi \\
\leq\left(m n R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{\frac{L_{1}}{R_{1}}-L_{2}}\right) e^{-L_{2}(1-x)}\|y\|_{B_{2}}^{2}
\end{array}
$$

provided that holds. A similar reasoning shows that

$$
\begin{aligned}
& \int_{t^{0}}^{t^{0}+T} \sum_{i=m+1}^{n}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} m_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}\left(s, \chi_{i}(s ; t, x)\right) d s\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \\
& \leq\left((n-m) n R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{2}+\frac{L_{1}}{R_{1}}}\right) e^{-L_{2}(1-x)}\|y\|_{B_{2}}^{2}
\end{aligned}
$$

4) For the remaining term, using a similar reasoning to the one used in the previous step, we obtain

$$
\begin{aligned}
& \int_{t^{0}}^{t^{0}+T} \sum_{i=1}^{n}\left|\sum_{j=1}^{n} \int_{\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right)}^{t} g_{i j}\left(s, \chi_{i}(s ; t, x)\right) y_{j}(s, 0) d s\right|^{2} e^{-L_{1}\left(t-t^{0}\right)} d t \\
& \leq\left(n^{2} R_{1} R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{1}}\right) e^{-L_{2}(1-x)}\|y\|_{B_{2}}^{2}
\end{aligned}
$$

In summary, we have established the following estimate (provided that (100) holds):

$$
\begin{align*}
&\left\|\mathcal{F}\left(y^{1}\right)-\mathcal{F}\left(y^{2}\right)\right\|_{B_{2}}^{2} \leq 3\left((n-m) m \alpha^{2} \frac{R_{1}}{\varepsilon} e^{\frac{R_{1}}{\varepsilon}}\right)\|y\|_{B_{1}}^{2} \\
&+3\left(m n \frac{R_{1}}{\varepsilon} R_{2}^{2} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{2}}+m n R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{\frac{L_{1}}{R_{1}}-L_{2}}+(n-m) n R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{2}+\frac{L_{1}}{R_{1}}}\right. \\
&\left.+n^{2} R_{1} R_{2}^{2} \frac{1}{\varepsilon^{2}} e^{\frac{R_{1}}{\varepsilon}} \frac{1}{L_{1}}\right)\|y\|_{B_{2}}^{2} \tag{103}
\end{align*}
$$

Consequently, we see from (101) and (103) that $\mathcal{F}$ indeed satisfies the contraction property (98)

Remark A.4. It can be shown that the broad solution is also the classical solution if the data of the system are smooth enough. It then follows by standard approximation arguments that the broad solution is also the so-called weak solution. We recall that the notion of weak solution for (91) is obtained by multiplying (91) by a smooth function and integrating by parts, that is, a function $y:\left(t^{0},+\infty\right) \times(0,1) \longrightarrow \mathbb{R}^{n}$ is a weak solution to (91) if $y \in C^{0}\left(\left[t^{0},+\infty\right) ; L^{2}(0,1)^{n}\right)$ and if it satisfies:

$$
\begin{align*}
& \int_{0}^{1} y\left(t^{0}+T, x\right) \cdot \varphi\left(t^{0}+T, x\right) d x- \int_{0}^{1} y^{0}(x) \cdot \varphi\left(t^{0}, x\right) d x \\
&+\int_{t^{0}}^{t^{0}+T} \int_{0}^{1} y(t, x) \cdot\left(-\frac{\partial \varphi}{\partial t}(t, x)-\Lambda(t, x) \frac{\partial \varphi}{\partial x}(t, x)-\left(\frac{\partial \Lambda}{\partial x}(t, x)+M(t, x)^{\operatorname{Tr}}\right) \varphi(t, x)\right) d x d t \\
&+\int_{t^{0}}^{t^{0}+T} \int_{0}^{1} y(t, \xi) \cdot F(t, \xi)^{\operatorname{Tr}} \Lambda_{--}(t, 1) \varphi_{-}(t, 1) d \xi d t=0 \tag{104}
\end{align*}
$$

for every $T>0$ and every $\varphi \in C^{1}\left(\left[t^{0}, t^{0}+T\right] \times[0,1]\right)^{n}$ such that, for every $t \in\left[t^{0}, t^{0}+T\right]$,

$$
\begin{gathered}
\varphi_{+}(t, 1)=0 \\
\varphi_{-}(t, 0)=-\Lambda_{--}(t, 0)^{-1}\left(Q(t)^{\operatorname{Tr}} \Lambda_{++}(t, 0) \varphi_{+}(t, 0)+\left(\operatorname{Id}_{\mathbb{R}^{m}} \quad Q(t)^{\operatorname{Tr}}\right) \int_{0}^{1} G(t, x)^{\operatorname{Tr}} \varphi(t, x) d x\right)
\end{gathered}
$$

In (104), we denoted by $A^{\mathrm{Tr}}$ the transpose of a matrix $A$ and $v_{1} \cdot v_{2}$ denotes the canonical scalar product between two vectors $v_{1}, v_{2}$ of $\mathbb{R}^{n}$.

## A. 3 Justification of the formal computations

In this section, we finally rigorously prove that the transformations that we used all along this paper are preserving broad solutions. We show how it works only for the Fredholm transformation of Section 2.3 (because it is simpler to present) but the reasoning is general and can be used for the Volterra transformation of Section 2.2 as well. More precisely, the goal of this section is to prove the following result:

Proposition A.5. Let $H_{--}=\left(h_{i j}\right)_{1 \leq i, j \leq m}$, where $h_{i j}$ is the solution to the differential equation (57). Let $t^{0} \geq 0$ be fixed. Let $z^{0} \in L^{2}(0,1)^{n}$ and let $z$ be the broad solution to

$$
\left\{\begin{array}{l}
\frac{\partial z}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial z}{\partial x}(t, x)=G^{3}(t, x) z(t, 0) \\
z_{-}(t, 1)=0, \quad z_{+}(t, 0)=Q^{1}(t) z_{-}(t, 0) \\
z\left(t^{0}, x\right)=z^{0}(x)
\end{array}\right.
$$

Then, the function $\gamma$ defined by the Fredholm transformation (48) is the broad solution to

$$
\left\{\begin{array}{l}
\frac{\partial \gamma}{\partial t}(t, x)+\Lambda(t, x) \frac{\partial \gamma}{\partial x}(t, x)=G^{2}(t, x) \gamma(t, 0) \\
\gamma_{-}(t, 1)=\int_{0}^{1} F^{2}(t, \xi) \gamma(t, \xi) d \xi, \quad \gamma_{+}(t, 0)=Q^{1}(t) \gamma_{-}(t, 0) \\
\gamma\left(t^{0}, x\right)=\gamma^{0}(x)
\end{array}\right.
$$

We recall that $H$ is given by $(53), F^{2}$ is the solution of $(50), Q^{1}$ is provided by Proposition 2.1, $G^{2}$ is provided by Proposition 2.4 and, finally, $G^{3}$ is given by (44).
Remark A.6. Obviously, we could use the explicit expression 60 - 58 of the solution $H$ to simplify the forthcoming arguments but we choose not to do so and to only use the differential equation (57) in order to give a general procedure that can also be used to justify the formal computations of Section 2.2 as well.

A similar result to Proposition A.5 can be found in [CN19, Proposition 3.5]. Here we propose a different and self-contained proof, based on the following characterization of broad solutions:

Lemma A.7. A function $y:\left(t^{0},+\infty\right) \times(0,1) \longrightarrow \mathbb{R}^{n}$ is the broad solution to (91) if, and only if, $y$ has the regularity (95) and, for every $i \in\{1, \ldots, n\}$, for a.e. $t>t^{0}$ and a.e. $x \in(0,1)$, the function $s \mapsto y_{i}\left(s, \chi_{i}(s ; t, x)\right)$ belongs to $H^{1}\left(\bar{s}_{i}^{\text {in }}\left(t^{0} ; t, x\right), s_{i}^{\text {out }}(t, x)\right)$ and it satisfies the $O D E$ (92).

The proof of Lemma A.7 is not difficult, it simply relies on the properties 10 and 17 .

## Proof of Proposition A.5.

1) The required regularity

$$
\gamma \in C^{0}\left(\left[t^{0}, t^{0}+T\right] ; L^{2}(0,1)^{n}\right) \cap C^{0}\left([0,1] ; L^{2}\left(t^{0}, t^{0}+T\right)^{n}\right), \quad \forall T>0
$$

is clear since $z$ also has this regularity and $(t, x) \mapsto \int_{0}^{1} H_{--}(t, x, \xi) z_{-}(t, \xi) d \xi$ is continuous (see e.g. Remark 2.9.
2) The initial condition in the ODE formulation

$$
\gamma_{i}\left(\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right), \chi_{i}\left(\bar{s}_{i}^{\mathrm{in}}\left(t^{0} ; t, x\right) ; t, x\right)\right)=b_{i}(\gamma)(t, x)
$$

is not difficult to check by using the boundary condition $z_{-}(t, 1)=0$ with the definition (50) of $F^{2}$ and Fubini's theorem (case $s_{i}^{\mathrm{in}}(t, x)>t^{0}$ and $i \in\{1, \ldots, m\}$ ), the condition $H_{--}(t, 0, \xi)=0$ (case $s_{i}^{\text {in }}(t, x)>t^{0}$ and $i \in\{m+1, \ldots, n\}$ ) and the definition of $\gamma^{0}$ (case $\left.s_{i}^{\text {in }}(t, x)<t^{0}\right)$.
3) It remains to check that, for every $i \in\{1, \ldots, n\}$, for a.e. $t>t^{0}$ and $x \in(0,1)$, the function $s \mapsto \gamma_{i}\left(s, \chi_{i}(s ; t, x)\right)$ belongs to $H^{1}\left(\bar{s}_{i}^{\text {in }}\left(t^{0} ; t, x\right), s_{i}^{\text {out }}(t, x)\right)$ with

$$
\begin{equation*}
\frac{d}{d s} \gamma_{i}\left(s, \chi_{i}(s ; t, x)\right)=\sum_{j=1}^{m} g_{i j}^{2}\left(s, \chi_{i}(s ; t, x)\right) \gamma_{j}(s, 0) \tag{105}
\end{equation*}
$$

By definition (48) of $\gamma$, we have

$$
\gamma_{i}\left(s, \chi_{i}(s ; t, x)\right)=z_{i}\left(s, \chi_{i}(s ; t, x)\right)-\sum_{j=1}^{m} \int_{0}^{1} h_{i j}\left(s, \chi_{i}(s ; t, x), \xi\right) z_{j}(s, \xi) d \xi
$$

For $i \in\{m+1, \ldots, n\}$, the identity 105 easily follows from the equation satisfied by $z_{i}$, the relation $z(\cdot, 0)=\gamma(\cdot, 0)$, and the fact that $h_{i j}=0$ for such indices (recall (53)).
Let us now assume that $i \in\{1, \ldots, m\}$. The equation satisfied by $z_{i}$ then gives

$$
\frac{d}{d s} z_{i}\left(s, \chi_{i}(s ; t, x)\right)=0, \quad \forall i \in\{1, \ldots, m\}
$$

On the other hand, since we know some information of $h_{i j}$ along the characteristic curve $s \mapsto$ $\chi_{i j}(s ; t, x, \theta)=\left(s, \chi_{i}(s ; t, x), \chi_{j}(s ; t, \theta)\right)$, we would like to perform the change of variable

$$
\xi=\chi_{j}(s ; t, \theta)
$$

where

$$
\begin{gathered}
a_{j}(s)=\theta_{j}(s ; t, 0), \quad b_{j}(s)=\theta_{j}(s ; t, 1), \\
\eta_{i j}(s, \theta)=h_{i j}\left(\chi_{i j}(s ; t, x, \theta)\right) z_{j}\left(s, \chi_{j}(s ; t, \theta)\right) \frac{\partial \chi_{j}}{\partial x}(s ; t, \theta) .
\end{gathered}
$$

We would like to use the formula

$$
\frac{d}{d s}\left(\int_{a_{j}(s)}^{b_{j}(s)} \eta_{i j}(s, \theta) d \theta\right)=b_{j}^{\prime}(s) \eta_{i j}\left(s, b_{j}(s)\right)-a_{j}^{\prime}(s) \eta_{i j}\left(s, a_{j}(s)\right)+\int_{a_{j}(s)}^{b_{j}(s)} \frac{\partial \eta_{i j}}{\partial s}(s, \theta) d \theta
$$

Clearly, $a_{j}, b_{j} \in C^{1}(\mathbb{R})$. Differentiating the relation $\xi=\chi_{j}\left(s ; t, \theta_{j}(s ; t, \xi)\right)$ with respect to $s$ we obtain

$$
a_{j}^{\prime}(s)=\frac{-\lambda_{j}(s, 0)}{\frac{\partial \chi_{j}}{\partial x}\left(s ; t, \theta_{j}(s ; t, 0)\right)}
$$

On the other hand, using (57) with $\xi=0$, 106) and the boundary condition $z_{-}(\cdot, 1)=0$, we have

$$
\eta_{i j}\left(s, a_{j}(s)\right)=-\frac{g_{i j}^{2}\left(s, \chi_{i}(s ; t, x)\right)}{\lambda_{j}(s, 0)} z_{j}(s, 0) \frac{\partial \chi_{j}}{\partial x}\left(s ; t, \theta_{j}(s ; t, 0)\right), \quad \eta_{i j}\left(s, b_{j}(s)\right)=0 .
$$

Using the ODEs satisfied along the characteristics by $h_{i j}$ (see (57)) and $z_{j}$, and using the relation (see (12))

$$
\frac{\partial^{2} \chi_{j}}{\partial s \partial x}(s ; t, \theta)=\frac{\partial \lambda_{j}}{\partial x}\left(s, \chi_{j}(s ; t, \theta)\right) \frac{\partial \chi_{j}}{\partial x}(s ; t, \theta)
$$

we can check that $\eta_{i j}$ has weak derivative with respect to $s$ which is equal to zero:

$$
\frac{\partial \eta_{i j}}{\partial s}(s, \theta)=0
$$

It follows from all the previous computations and the relation $z(\cdot, 0)=\gamma(\cdot, 0)$ that

$$
\frac{d}{d s} \gamma_{i}\left(s, \chi_{i}(s ; t, x)\right)=\sum_{j=1}^{m} a_{j}^{\prime}(s) \eta_{i j}\left(s, a_{j}(s)\right)=\sum_{j=1}^{m} g_{i j}^{2}\left(s, \chi_{i}(s ; t, x)\right) \gamma_{j}(s, 0)
$$

Thanks to 13 and the implicit function theorem there exists $\theta_{j} \in C^{1}\left(\mathbb{R}^{3}\right)$ such that, for every $(s, t, \xi) \in \mathbb{R}^{3}$, we have

$$
\begin{equation*}
\xi=\chi_{j}\left(s ; t, \theta_{j}(s ; t, \xi)\right), \quad \frac{\partial \theta_{j}}{\partial \xi}(s ; t, \xi)>0 \tag{106}
\end{equation*}
$$

Using this change of variable, we have

$$
\gamma_{i}\left(s, \chi_{i}(s ; t, x)\right)=z_{i}\left(s, \chi_{i}(s ; t, x)\right)-\sum_{j=1}^{m} \int_{a_{j}(s)}^{b_{j}(s)} \eta_{i j}(s, \theta) d \theta
$$

## B Constructions of $s_{i j}^{\text {in }}, s_{i j}^{\text {out }}$ and $s_{i j \ell}^{\text {disc }}$

In this appendix, we give a proof of the existence $s_{i j}^{\mathrm{in}}, s_{i j}^{\text {out }}$ and $s_{i j \ell}^{\text {disc }}$ satisfying the properties stated in Proposition 3.3 and Proposition 3.5 . We will make use of the following simple lemma:

Lemma B.1. Let $f \in C^{1}([a, b])(a<b)$ satisfy the following property:

$$
\begin{equation*}
\forall s \in[a, b], \quad f(s)=0 \quad \Longrightarrow \quad f^{\prime}(s)<0 \tag{107}
\end{equation*}
$$

Then, there exists a unique $c \in[a, b]$ such that

$$
f(s)>0, \quad \forall s \in(a, c), \quad f(s)<0, \quad \forall s \in(c, b)
$$

Moreover, $c$ has the properties listed in Table 1 (an Ø means that such a situation can not occur).

|  | $f(b)>0$ | $f(b)=0$ | $f(b)<0$ |
| :---: | :---: | :---: | :---: |
| $f(a)>0$ | $c=b$ | $c=b$ | $f(c)=0$ |
| $f(a)=0$ | $\emptyset$ | $\emptyset$ | $c=a$ |
| $f(a)<0$ | $\emptyset$ | $\emptyset$ | $c=a$ |

Table 1: Properties of $c$

Proof of Proposition 3.3. We recall that $i \in\{1, \ldots, m\}$ and we refer to Figures 2, 3, 4 and 5 for a clarification of the geometric situation (at a fixed $t$ ). We only focus on the existence part since the uniqueness readily follows from the properties that have to be satisfied.

1) Assume that $j \in\{1, \ldots, i-1\}$. For every $(t, x, \xi) \in \overline{\mathcal{P}}$ such that $x<1$, we introduce the $C^{1}$ function

$$
f: s \in\left[\max \left\{s_{i}^{\mathrm{in}}(t, x), s_{j}^{\mathrm{in}}(t, \xi)\right\}, t\right] \mapsto \chi_{j}(s ; t, \xi)-\chi_{i}(s ; t, x)
$$

Note that the interval has a non empty interior since $x<1$ and $\xi \leq x<1$ (see (15)- 16 ). This function clearly satisfies the property (107) thanks to the ODE (9) and the assumption (3) since $j<i$. Consequently, Lemma B. 1 applies and gives the existence of $s_{i j}^{\mathrm{in}}(t, x, \xi)$ with

$$
\max \left\{s_{i}^{\mathrm{in}}(t, x), s_{j}^{\mathrm{in}}(t, \xi)\right\} \leq s_{i j}^{\mathrm{in}}(t, x, \xi) \leq t
$$

and such that

$$
\chi_{j}(s ; t, \xi)<\chi_{i}(s ; t, x), \quad \forall s \in\left(s_{i j}^{\mathrm{in}}(t, x, \xi), t\right)
$$

Clearly, $(t, x, \xi) \mapsto t-s_{i j}^{\text {in }}(t, x, \xi) \in L^{\infty}(\mathcal{P})$ thanks to 19 . Moreover, it follows from Table 1 that

$$
\begin{cases}s_{i j}^{\text {in }}(t, x, \xi)=t & \text { if } s_{i}^{\text {in }}(t, x)<s_{j}^{\text {in }}(t, \xi) \text { and } \xi=x  \tag{108}\\ f\left(s_{i j}^{\text {in }}(t, x, \xi)\right)=0 & \text { if } s_{i}^{\text {in }}(t, x)<s_{j}^{\text {in }}(t, \xi) \text { and } \xi<x \\ s_{i j}^{\text {in }}(t, x, \xi)=s_{i}^{\text {in }}(t, x) & \text { if } s_{i}^{\text {in }}(t, x)=s_{j}^{\text {in }}(t, \xi) \text { and } \xi<x, \\ s_{i j}^{\text {in }}(t, x, \xi)=s_{i}^{\text {in }}(t, x) & \text { if } s_{i}^{\text {in }}(t, x)>s_{j}^{\text {in }}(t, \xi) \text { and } \xi<x\end{cases}
$$

Let us now complete the definition of $s_{i j}^{\text {in }}$ on the remaining parts of $\overline{\mathcal{P}}$. The missing case in (108) is when $s_{i}^{\text {in }}(t, x) \geq s_{j}^{\mathrm{in}}(t, \xi)$ and $\xi=x$. However, unless $x=1$, these conditions are not compatible since $s_{i}^{\text {in }}(t, x)<s_{j}^{\text {in }}(t, x)$ for $j<i \leq m$ (see 21). Consequently, it only remains to define $s_{i j}^{\mathrm{in}}$ in the part where $x=1$, which we do now by setting

$$
\begin{equation*}
s_{i j}^{\mathrm{in}}(t, 1, \xi)=t \tag{109}
\end{equation*}
$$

We can check that $s_{i j}^{\text {in }}$ defined by $(108)-109$ belongs to $C^{0}(\overline{\mathcal{P}})$ (for the second case in 108 this follows from the implicit function theorem). Therefore, such a $s_{i j}^{\text {in }}$ clearly satisfies all the properties claimed in the statement of item (i) of Proposition 3.3
2) Assume that $j=i$. We will show that, in this case, we can simply take

$$
s_{i i}^{\text {out }}(t, x, \xi)=s_{i}^{\text {out }}(t, \xi)
$$

Clearly, $s_{i i}^{\text {out }} \in C^{0}(\overline{\mathcal{P}})$ with $(t, x, \xi) \mapsto s_{i i}^{\text {out }}(t, x, \xi)-t \in L^{\infty}(\mathcal{P})$ thanks to (19) and $s_{i i}^{\text {out }}(t, x, \xi)>t$ as long as $\xi>0$ (see 15 - 16 ). Let us now observe that, for $\xi<x$, we have from (13):

$$
\chi_{i}(s ; t, \xi)<\chi_{i}(s ; t, x), \quad \forall s \in \mathbb{R}
$$

and $\chi_{i}(s ; t, \xi)>0$ for $s \in\left(t, s_{i i}^{\text {out }}(t, x, \xi)\right)$ since $s_{i}^{\text {out }}(t, \xi)<s_{i}^{\text {out }}(t, x)$ by (recall that $j=i \in$ $\{1, \ldots, m\})$.
3) The proof for the case $j \in\{i+1, \ldots, m\}$ is similar to the proof of part 1) by considering, for each $(t, x, \xi) \in \overline{\mathcal{P}}$ such that $\xi, x>0$, the function

$$
\begin{equation*}
f: s \in\left[t, \min \left\{s_{i}^{\text {out }}(t, x), s_{j}^{\text {out }}(t, \xi)\right\}\right] \mapsto \chi_{i}(s ; t, x)-\chi_{j}(s ; t, \xi) \tag{110}
\end{equation*}
$$

4) The proof for the case $j \in\{m+1, \ldots, n\}$ is also similar to the proof of part 1) by considering, for each $(t, x, \xi) \in \overline{\mathcal{P}}$ such that $0 \leq \xi<x \leq 1$, the function $f$ defined again by 110 .

Proof of Proposition 3.5. The difference with the proof of Proposition 3.3 is that we do not need to neither track the regularity of the point where the function $f$ vanishes nor its sign on the left and right of this zero. It is a straightforward consequence of Lemma B. 1 applied to the following functions (it is enough to consider non empty intervals):

1) For $j \in\{1, \ldots, i-1\}$, we use

$$
f: s \in\left[s_{i j}^{\mathrm{in}}(t, x, \xi), t\right] \mapsto \chi_{j}(s ; t, \xi)-\psi_{i \ell}\left(s, \chi_{i}(s ; t, x)\right)
$$

Using the ODE (9) satisfied by $\chi_{j}$ and using the equation satisfied by $\psi_{i \ell}$, we have

$$
f^{\prime}(s)=\lambda_{j}\left(s, \chi_{j}(s ; t, \xi)\right)-\lambda_{\ell}\left(s, \psi_{i \ell}\left(s, \chi_{i}(s ; t, x)\right)\right)
$$

Since $j<\ell$, this shows that such a $f$ satisfies the property (107) of Lemma B. 1
2) For $j \in\{i, \ldots, \ell-1\}$, we use

$$
\begin{equation*}
f: s \in\left[t, s_{i j}^{\text {out }}(t, x, \xi)\right] \mapsto \chi_{j}(s ; t, \xi)-\psi_{i \ell}\left(s, \chi_{i}(s ; t, x)\right) \tag{111}
\end{equation*}
$$

3) For $j \in\{\ell+1, \ldots, m\}$, we use the function $-f$, where $f$ is given by 111 .
4) For $j \in\{m+1, \ldots, n\}$, we use the same function $f$ given by (in fact, the result then directly follows from the intermediate value theorem).

## C Construction of $\Omega_{i}$

This appendix is devoted to the proof of Lemma 3.7, that is to the existence of the key change of variable needed in the proof of Proposition 3.6. We recall that $i \in\{1, \ldots, m\}$.

1) Inspired by the time-independent case (see Remark 3.6), we look for $\Omega_{i}$ in the following form:

$$
\Omega_{i}(t, x, \xi)=\omega_{i}^{1}(t, x)-\omega_{i}^{\nu}(t, \xi)
$$

where, at each fixed $\nu \in(0,1], \omega_{i}^{\nu}(\cdot, \cdot)$ is the solution to the following linear hyperbolic equation:

$$
\left\{\begin{array}{l}
\frac{\partial \omega_{i}^{\nu}}{\partial t}(t, x)+\frac{\lambda_{i}(t, x)}{\nu} \frac{\partial \omega_{i}^{\nu}}{\partial x}(t, x)=0,  \tag{112}\\
\omega_{i}^{\nu}(t, 0)=t
\end{array} \quad t \in \mathbb{R}, x \in[0,1]\right.
$$

The solution of 112 is explicit:

$$
\begin{equation*}
\omega_{i}^{\nu}(t, x)=\omega_{i}^{\nu}\left(s_{i}^{\text {out }, \nu}(t, x), 0\right)=s_{i}^{\text {out }, \nu}(t, x) \tag{113}
\end{equation*}
$$

where $s_{i}^{\text {out, } \nu}(t, x) \geq t$ (with $s_{i}^{\text {out, } \nu}(t, x)=t \Longleftrightarrow x=0$ ) is the unique number such that

$$
\begin{equation*}
\chi_{i}^{\nu}\left(s_{i}^{\text {out }, \nu}(t, x) ; t, x\right)=0, \tag{114}
\end{equation*}
$$

where $s \mapsto \chi_{i}^{\nu}(s ; t, x)$ is the solution to the ODE

$$
\left\{\begin{array}{l}
\frac{\partial \chi_{i}^{\nu}}{\partial s}(s ; t, x)=\frac{1}{\nu} \lambda_{i}\left(s, \chi_{i}^{\nu}(s ; t, x)\right), \quad \forall s \in \mathbb{R}  \tag{115}\\
\chi_{i}^{\nu}(t ; t, x)=x
\end{array}\right.
$$

We can check that the map $(t, x, \nu) \mapsto \omega_{i}^{\nu}(t, x)$ belongs to $C^{1}(\mathbb{R} \times[0,1] \times(0,1])$.
2) We now prove that there exists $\delta>0$ such that, for every $t \in \mathbb{R}, x \in[0,1]$ and $\nu \in(0,1]$,

$$
\begin{equation*}
\frac{\partial \omega_{i}^{\nu}}{\partial t}(t, x) \geq \varepsilon \delta, \quad \frac{\partial \omega_{i}^{\nu}}{\partial x}(t, x) \geq \nu \delta, \quad \frac{\partial \omega_{i}^{\nu}}{\partial \nu}(t, x) \geq 0 \tag{116}
\end{equation*}
$$

Using the equation (112) and the assumption (3), it is clear that the estimate for $\partial \omega_{i}^{\nu} / \partial t$ follows from the estimate of $\partial \omega_{i}^{\nu} / \partial x$. Note from (113) that $\partial \omega_{i}^{\nu} / \partial x=\partial s_{i}^{\text {out, } \nu} / \partial x$. Taking the derivative of 114 with respect to $x$, we obtain

$$
\frac{1}{\nu} \lambda_{i}\left(s_{i}^{\text {out }, \nu}(t, x), \chi_{i}^{\nu}\left(s_{i}^{\text {out }, \nu}(t, x) ; t, x\right)\right) \frac{\partial s_{i}^{\text {out }, \nu}}{\partial x}(t, x)+\frac{\partial \chi_{i}^{\nu}}{\partial x}\left(s_{i}^{\text {out }, \nu}(t, x) ; t, x\right)=0 .
$$

Since $\lambda_{i} \in L^{\infty}(\mathbb{R} \times(0,1))$, we have to bound $\frac{\partial \chi_{i}^{\nu}}{\partial x}\left(s_{i}^{\text {out }, \nu}(t, x) ; t, x\right)$ from below by a positive constant that does not depend on $t, x$ and $\nu$. From 115 we can show that

$$
\frac{\partial \chi_{i}^{\nu}}{\partial x}(s ; t, x)=e^{\frac{1}{\nu} \int_{t}^{s} \frac{\partial \lambda_{i}}{\partial x}\left(\theta, \chi_{i}^{\nu}(\theta ; t, x)\right) d \theta}
$$

so that

$$
\frac{\partial \chi_{i}^{\nu}}{\partial x}\left(s_{i}^{\text {out }, \nu}(t, x) ; t, x\right) \geq e^{\frac{1}{\nu}\left(s_{i}^{\text {out }, \nu}(t, x)-t\right) \inf _{\mathbb{R} \times[0,1]} \frac{\partial \lambda_{i}}{\partial x}}
$$

This establishes the desired lower bound since $\frac{\partial \lambda_{i}}{\partial x} \in L^{\infty}(\mathbb{R} \times(0,1))$ and $0 \leq s_{i}^{\text {out }, \nu}(t, x)-t \leq \frac{x}{\varepsilon} \nu$ (the proof is similar to the one of $(19)$ ). Note that it follows as well from this estimate that $\Omega_{i} \in L^{\infty}(\mathbb{R} \times(0,1) \times(0,1))$.

To prove the remaining estimate in (116), we denote by $\gamma^{\nu}=\partial \omega_{i}^{\nu} / \partial \nu$ and observe that it satisfies

$$
\left\{\begin{array}{l}
\frac{\partial \gamma^{\nu}}{\partial t}(t, x)+\frac{\lambda_{i}(t, x)}{\nu} \frac{\partial \gamma^{\nu}}{\partial x}(t, x)=\frac{\lambda_{i}(t, x)}{\nu^{2}} \frac{\partial \omega_{i}^{\nu}}{\partial x}(t, x) \leq 0, \\
\gamma^{\nu}(t, 0)=0
\end{array} \quad t \in \mathbb{R}, x \in[0,1]\right.
$$

It immediately follows that $\gamma^{\nu} \geq 0$.
3) Let us now check the estimates (87) and 88). We have

$$
\left.\begin{array}{rl}
\frac{\partial \Omega_{i}}{\partial t}(t, x, \nu)+ & \lambda_{i}(t, x) \frac{\partial \Omega_{i}}{\partial x}(t, x, \nu)+\lambda_{j}(t, \xi) \frac{\partial \Omega_{i}}{\partial \xi}(t, x, \nu) \\
& =\frac{\partial \omega_{i}^{1}}{\partial t}(t, x)+\lambda_{i}(t, x) \frac{\partial \omega_{i}^{1}}{\partial x}(t, x)-\frac{\partial \omega_{i}^{\nu}}{\partial t}(t, \xi)
\end{array}\right) \quad \lambda_{j}(t, \xi) \frac{\partial \omega_{i}^{\nu}}{\partial x}(t, \xi) .
$$

Since $\lambda_{j} / \lambda_{i} \leq 1$ for $i \leq j$ and $i \leq m$, we see that the estimate 88 is obtained by simply taking

$$
\begin{equation*}
0<\varepsilon_{0} \leq \varepsilon \delta(1-\nu), \quad 0<\nu<1 \tag{117}
\end{equation*}
$$

On the other hand, let us introduce

$$
r=\max _{1 \leq j<i} \sup _{\substack{t \in \mathbb{R} \\ \xi \in[0,1]}} \frac{\lambda_{i}(t, \xi)}{\lambda_{j}(t, \xi)} .
$$

Clearly, $0<r \leq 1$. In fact, $r<1$ since from we have, for $j<i \leq m$,

$$
\sup _{\substack{t \in \mathbb{R} \\ \xi \in[0,1]}} \frac{\lambda_{i}(t, \xi)}{\lambda_{j}(t, \xi)} \leq 1-\frac{\varepsilon}{\left\|\lambda_{j}\right\|_{L^{\infty}}}
$$

The estimate (87) now follows from (116) by taking

$$
\begin{equation*}
0<\varepsilon_{0} \leq \varepsilon \delta\left(\frac{\nu}{r}-1\right), \quad r<\nu \leq 1 \tag{118}
\end{equation*}
$$

Note that the conditions (117) and (118) are compatible by taking $\nu$ close enough to 1.
4) It remains to check that $\Omega_{i} \geq 0$ on $\overline{\mathcal{P}}$. Since both functions $\nu \mapsto \omega_{i}^{\nu}(t, x)$ and $x \mapsto \omega_{i}^{\nu}(t, x)$ are nondecreasing by 116 and $\xi \leq x$, we have

$$
\omega_{i}^{1}(t, x) \geq \omega_{i}^{\nu}(t, x) \geq \omega_{i}^{\nu}(t, \xi)
$$

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