# NUMERICAL METHOD FOR THE COMPUTATION OF TANGENT VECTORS TO $2 \times 2$ HYPERBOLIC SYSTEMS OF CONSERVATION LAWS 

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#### Abstract

We are interested in the development of a numerical method for solving optimal control problems governed by hyperbolic systems of conservation laws. The main difficulty of computing the derivative in the case of shock waves is resolved in the presented scheme. Our approach is based on a combination of a relaxation approach in combination with a numerical scheme to resolve the evolution of the tangent vectors. Numerical results for optimal control problems are presented.


Key words. Conservation laws, optimization, tangent vectors.
AMS subject classifications. 35L65, 49K20, 49 K 40

1. Introduction We are concerned with a numerical approach to optimization problems governed by systems of hyperbolic partial differential equations in a single spatial dimension. As a prototype, we consider a tracking type problem for a terminal state $y_{d}$ prescribed at some given time $t=T$ and the control acts as initial condition $u_{0}$. A mathematical formulation of this optimal control problem is reduced to minimizing a functional, and, for instance, it can be stated as follows:

$$
\begin{equation*}
\min _{u_{0}} J\left(y, y_{d}\right), \tag{1.1}
\end{equation*}
$$

where $J$ is the given cost functional and $y \in \mathbb{R}^{n}$ is the unique entropy solution of the nonlinear conservation law

$$
\begin{align*}
y_{t}+f(y)_{x} & =0, & & x \in \mathbb{R}, \quad t>0, \\
y(0, x) & =u_{0}(x), & & x \in \mathbb{R} . \tag{1.2}
\end{align*}
$$

There has been tremendous progress in both analytical and numerical studies of problems of type (1.1), (1.2), see, e.g., [1, 2, 3, 7, 27, 12, 8, 13, 19, 24, 29, 28, 20, 18, 22]. Its solution relies on the property of the evolution operator $\mathcal{S}_{t}: u_{0}(\cdot) \rightarrow y(\cdot, t)=\mathcal{S}_{t} u_{0}(\cdot)$ for (1.2). It is known that the semi-group $\mathcal{S}_{t}$ generated by a nonlinear hyperbolic conservation law is generically nondifferentiable in $L^{1}$ even in the scalar one-dimensional (1-D) case (see, e.g., [12, Example 1]). A calculus for the first-order variations of $\mathcal{S}_{t} u_{0}$ with respect to $u_{0}$ has been established in [12, Theorems 2.2 and 2.3] for general 1-D systems of conservation laws with a piecewise Lipschitz continuous $u_{0}$ that contains finitely many discontinuities. Therein, the concept of generalized first order tangent vectors has been introduced to characterize the evolution of variations with respect to $u_{0}$, see [12, equations (2.16)-(2.18)]. It has been further extended in [11] to establish continuous dependence of $\mathcal{S}_{t}$ on the initial data $u_{0}$. This result has been extended to BV initial data in $[7,3]$ and lead to the introduction of a differential structure for $u_{0} \rightarrow \mathcal{S}_{t} u_{0}$, called shift-differentiability, see e.g. [3, Definition 5.1]. Further extensions have also been discussed for example [17]. Related to that equations for the generalized cotangent vectors have been introduced for 1-D systems in [9, Proposition 4]. These

[^0]equations (also called adjoint equations) consist of a nonconservative transport equation [9, equation (4.2)] and an ordinary differential equation [9, equations (4.3)-(4.5)] for the tangent vector and shift in the positions of possible shocks in $y(x, t)$, respectively. Necessary conditions for a general optimal control problem have been established in [9, Theorem 1]. However, this result was obtained using strong assumptions on $u_{0}$ (see [9, Remark 4] and [3, Example 5.5]), which in the 1-D scalar case can be relaxed as shown for example in $[29,13]$. We note that the nonconservative transport part of the adjoint equation has been intensively studied also independently from the optimal control context. In the scalar case we refer to $[29,5,6,26]$ for a notion of solutions and properties of solutions to those equations. Analytical results for optimal control problems in the case of a scalar hyperbolic conservation law with a convex flux have also been developed using a different approach in [29]. The relation to the weak formulation has been discussed in [2] in the case of Burgers' equation.

Numerical methods for the optimal control problems have been discussed in [1, 20]. In [19], the adjoint equation has been discretized using a Lax-Friedrichs-type scheme, obtained by including conditions along shocks and modifying the Lax-Friedrichs numerical viscosity. Convergence of the modified Lax-Friedrichs scheme has been rigorously proved in the case of a smooth convex flux function. Convergence results have also been obtained in [29] for the class of schemes satisfying the one-sided Lipschitz condition (OSLC) and in [1] for implicit-explicit finite-volume methods. Other examples of finite volume methods and Lagrangian methods are given in [14, 23].

In [13], analytical and numerical results for the optimal control problem (1.1) coupled with the 1-D inviscid Burgers' equation have been presented in the particular case of a least-square cost functional $J$. Therein, existence of a minimizer $u_{0}$ was proven, however, uniqueness could not be obtained for discontinuous functions $y$. This result was also extended to the discretized optimization problem provided that the numerical schemes satisfy either the OSLC or discrete Oleinik's entropy condition. Furthermore, convergence of numerical schemes was investigated in the case of convex flux functions and with a-priori known shock positions, and numerical resolution of the adjoint equations in both the smooth and non-smooth cases was studied. In [21] perturbations of initial data are studied using an additional spatial dimension. Numerical results as well as a formalism to derive the linearized equations have been presented therein. In the scalar case of a production model coupled to ordinary differential equations has been studied in [16]. Therein, convergence of the wave-front tracking approximation to the tangent vector equation is proven.

We contribute to the discussion by introducing a novel scheme which allows to include the arising discontinuities in an optimization framework and without an a-priori assumption on the location of the discontinuities. This is possible and computationally efficient under three basic assumptions: first, we only compute derivatives with respect to piecewise constant controls $u_{0}$; second, the system (1.2) is not solved directly, but an $\epsilon$-relaxation approximation (2.2) is solved instead and last, we compute the exact derivative for the $\epsilon$-approximation of the system (2.2) using tangent vectors. The number of discontinuities in $u_{0}$ may herein be as large as $\frac{1}{\Delta x}$, where $\Delta x$ is the spatial width of the numerical grid. The overall algorithm requires the solution of two additional hyperbolic partial differential equations. The motivation and theoretical investigations are presented in section 2 and numerical results in section 3. The results presented are on the relaxation formulation of the Burgers' equation as well as linear hyperbolic systems. Nevertheless, the approach also applies to nonlinear hyperbolic systems. In this case the current approach leads to higher diffusion due to the requirement of a single
characteristic speed bounding all characteristic speed as seen in the implementation for example in equation (2.25).

## 2. Motivation and theoretical results

In order to derive a numerical scheme we consider a relaxation approximation [25] to (1.2). For simplicity we consider only the Jin-Xin relaxation in the case $n=1$ and on the full real line $x \in \mathbb{R}$. Then, the hyperbolic relaxation for

$$
\begin{equation*}
y_{t}^{(1)}+f\left(y^{(1)}\right)_{x}=0, y^{(1)}(0, x)=u_{0}(x), \tag{2.1}
\end{equation*}
$$

is given by

$$
y_{t}+\left(\begin{array}{cc}
0 & 1  \tag{2.2}\\
a^{2} & 0
\end{array}\right) y_{x}=\binom{0}{\frac{1}{\epsilon}\left(f\left(y^{(1)}\right)-y^{(2)}\right)}
$$

and by the initial data

$$
\begin{equation*}
y^{(1)}(0, x)=u_{0}(x), y^{(2)}(0, x)=f\left(u_{0}(x)\right) . \tag{2.3}
\end{equation*}
$$

We assume that the value $a^{2}$ fulfills the sub characteristic condition

$$
\begin{equation*}
a^{2}-f^{\prime}\left(y^{(1)}\right)^{2} \geq 0 \forall y^{(1)} \tag{2.4}
\end{equation*}
$$

For positive $\epsilon$ it is known that a solution $y$ yields an approximation to (1.2) in the following sense. For $\epsilon$ sufficiently small we have up to second order in $\epsilon$

$$
y_{t}^{(1)}+f\left(y^{(1)}\right)_{x}=\epsilon\left(\left(a^{2}-f^{\prime}\left(y^{(1)}\right)^{2}\right) y_{x}^{(1)}\right)_{x}
$$

Formally we obtain the original conservation law for $y^{(1)}$ in the limit $\epsilon=0$. We refer to [4] for a detailed analysis. For $\epsilon>0$ the method is known as a relaxing scheme. The main advantage of (2.2) over (1.2) is the linear transport which greatly simplifies the computation of associated tangent vectors (in particular the numerical resolution of the shock variations will turn out trivial for (2.2)).

In the following we will therefore discuss the optimization problem (1.1) with respect to (2.2). Still, the numerical computation of tangent vectors in general poses severe challenges addressed below. Therefore, we further simplify by only considering piecewise constant controls $u_{0}$. In the following $T V(\cdot)$ denotes the total variation.

Given some $C>0$, we consider problem (1.1) subject to (2.2) and (2.3) for controls $u_{0} \in \mathcal{U}$.
Definition 2.1. We indicate by $\mathcal{U}:=\{u: \mathbb{R} \rightarrow \mathbb{R}: u$ measurable,$T V(u) \leq C$, $u$ piecewise constant $\}$ the set of admissible controls. For every $u \in \mathcal{U}$ we indicate by $x_{k}=x_{k}(u), k=1, \ldots, N(u)$ the points of discontinuity of $u$.

For $y_{d} \in L^{1}(\mathbb{R})$, some $T>0$ and a bounded interval $I \subset \mathbb{R}$ we consider as a prototype example an unregularized cost functional of tracking type

$$
\begin{equation*}
J\left(y, y_{d} ; T, I\right)=\int \chi_{I}(x)\left(y^{(1)}(T, x)-y_{d}(x)\right)^{2} d x \tag{2.5}
\end{equation*}
$$

The dependence of $J$ on terminal time $T$ and the interval $I$ will be dropped from now on. We assume that both are fixed and will not change throughout the manuscript.

We now introduce the notion of tangent vectors, see [12] and [11]. In (2.2) we assume $\epsilon>0$ fixed, $f \in C^{4}(\mathbb{R})$ and $a^{2}$ fulfills the subcharacteristic condition. For a function $u \in \mathcal{U}$


FIg. 2.1. Construction of a tangent vector $\left(\xi_{i}, v\right)$ to $u_{\delta}$ having a discontinuity at $x_{i}$.
a generalized tangent vector consists of two components $(v, \xi)$ where $v \in L^{1}(\mathbb{R})$ describes the $L^{1}$ infinitesimal displacement and $\xi \in \mathbb{R}^{N(u)}$ describes the infinitesimal displacement of $N(u)$ discontinuities. A norm on the space of tangent vectors $T_{u}:=L^{1}\left(\mathbb{R} ; \mathbb{R}^{n}\right) \times \mathbb{R}^{N(u)}$ is given by

$$
\begin{equation*}
\|(v, \xi)\|:=\|v\|_{L^{1}}+\sum_{i=1}^{N(u)}\left|\Delta_{i} u\right|\left|\xi_{i}\right| \tag{2.6}
\end{equation*}
$$

where $\Delta_{i} u=u\left(x_{i}+\right)-u\left(x_{i}-\right)$ denotes the jump in $u$. The norm depends on $u$ through the number of points of discontinuity. Tangent vectors may be used to describe variations of $u$ as exemplified in Figure 2.1. For $\delta>0$ an infinitesimal displacement $u_{\delta}$ of $u$ is given by

$$
\begin{equation*}
u_{\delta}=u+\delta v-\sum_{i: \xi_{i}>0}^{N(u)} \Delta_{i} u \chi_{\left[x_{i}, x_{i}+\delta \xi_{i}\right]}+\sum_{i: \xi_{i}<0}^{N(u)} \Delta_{i} u \chi_{\left[x_{i}+\delta \xi_{i}, x_{i}\right]} . \tag{2.7}
\end{equation*}
$$

$u_{\delta}$ is obtained from $u$ by shifting the function values by $\delta v$ and the $i$ th discontinuity by $\delta \xi_{i}$. For $\delta$ sufficiently small $u_{\delta}$ has the same number of discontinuities as $u$. Note that if $\xi \neq 0$ then the function $\delta \rightarrow u_{\delta}$ is not differentiable in $L^{1}$ as the ratio $\frac{u_{\delta+h}-u_{\delta}}{h}$ does not converge to any limit in $L^{1}$ for $h \rightarrow 0$. However, the previous limit remains meaningful if interpreted as a weak limit in a space of measures with a singular point mass located at $x_{i}$ and having magnitude $\left|\Delta_{i} u\right| \xi_{i}$. Therefore, in [12] a class of variations $\delta \rightarrow u_{\delta}$ is described up to first order by (generalized) tangent vectors $(v, \xi)$.

In general, the norm (2.6) depends on the number of discontinuities $N(u)$ and therefore on $u$. As a numerical method we will restrict ourselves to piecewise constant controls $u_{0}$. We are going study the number $N\left(u_{0}\right)$ equal to the number of grid cells in the spatial domain. This number is known a priori.

We introduce some notation and definitions (see [12, 11]) already in view of the special system (2.2). Let $u \in L^{1}\left(\mathbb{R} ; \mathbb{R}^{n}\right)$ be a piecewise Lipschitz continous function with $N=N(u)$ jumps. Consider $\Sigma_{u}$ the family of all continuous paths $\gamma:\left[0, \delta_{0}\right] \rightarrow L_{\text {loc }}^{1}$
with $\gamma(0)=u$ with $\delta_{0}$ possibly depending on $\gamma$. We recall the following definition [12, Definition 1,3].
Definition 2.2. The space of generalized tangent vectors to a piecewise Lipschitz function $u$ with jumps at the points $x_{1}<x_{2} \cdots<x_{N}$ is $T_{u}:=L^{1}\left(\mathbb{R} ; \mathbb{R}^{n}\right) \times \mathbb{R}^{N}$. A continuous path $\gamma \in \Sigma_{u}$ generates a tangent vector $(v, \xi) \in T_{u}$ if

$$
\lim _{\delta \rightarrow 0} \frac{1}{\delta}\|\gamma(\delta)-\bar{\gamma}(\delta)\|_{L^{1}}=0
$$

for

$$
\bar{\gamma}:=u+\delta v-\sum_{i: \xi_{i}>0} \Delta_{i} u \chi_{\left[x_{i}, x_{i}+\delta \xi_{i}\right]}+\sum_{i: \xi_{i}<0} \Delta_{i} u \chi_{\left[x_{i}+\delta \xi_{i}, x_{i}\right]} .
$$

Let u be a piecewise Lipschitz function with simple discontinuities [12, Definition 2]. Then, a path $\gamma \in \Sigma_{u}$ is a regular variation for $u$ if additionally all function $\gamma(\delta)=u_{\delta}$ are piecewise Lipschitz with simple discontinuities and the jumps $x_{i}^{\delta}$ depend continuously on $\delta$. A regular variation $\gamma$ for $u$ generates a tangent vector $(v, \xi)$ by

$$
\begin{equation*}
\xi_{i}=\lim _{\delta \rightarrow 0} \frac{x_{i}^{\delta}-x_{i}}{\delta}, \lim _{\delta \rightarrow 0} \int_{a}^{b}\left\|\frac{u^{\delta}\left(x_{i}^{\delta}+y\right)-u\left(x_{i}+y\right)}{\delta}-v\left(x_{i}+y\right)-\xi_{i} u_{x}\left(x_{i}+y\right)\right\| d y=0 \tag{2.8}
\end{equation*}
$$

whenever $\left[x_{i}+a, x_{i}+b\right]$ does not contain any other point of discontinuity of $u$ except $x_{i}$. Further, the length of a regular path $\gamma$ can be computed by (2.6). We now consider the initial data $u_{0}$ and a regular variation generating the tangent vector $(v, \xi) \in T_{u}$. Under suitable regularity assumptions [12, Theorem 2.2] regular variations are locally preserved by the system (2.2) and linearized equations for the time evolution of the tangent vector $(v(t, \cdot), \xi(\cdot))$ can be derived. We have the following result [12, Theorem 2.2].

Lemma 2.1. Let $y(\cdot, \cdot)$ be a piecewise Lipschitz continuous solution to (2.2) and initial data (2.3) $y(0, \cdot)=\bar{y}$ piecewise Lipschitz with $N$ simple discontinuities. Let $(\bar{v}, \bar{\xi}) \in T_{\bar{y}}$ be a tangent vector to $\bar{y}$ generated by the regular variation $\gamma$ with $\gamma(\delta)=\bar{y}_{\delta}$. Let $y_{\delta}(t, x)$ be the solution to (2.2) and initial data (2.3) $y_{\delta}(0, x)=\bar{y}_{\delta}(\cdot)$. Then, there exists a time $t_{0}>0$ such that for all $t \in\left[0, t_{0}\right]$ the path $\bar{\gamma}$ with $\bar{\gamma}(\delta)=y_{\delta}(t, \cdot)$ is a regular variation of $y(t, \cdot)$ generating the tangent vector $(v(t), \xi(t)) \in T_{y(t, \cdot)}$. Further, $(v, \xi)$ is the unique broad solution to

$$
v(0, \cdot)=\bar{v}(\cdot), \quad v_{t}+\left(\begin{array}{rr}
0 & 1  \tag{2.9}\\
a^{2} & 0
\end{array}\right) v_{x}=\frac{1}{\epsilon}\binom{0}{f^{\prime}\left(y^{(1)}\right) v^{(1)}-v^{(2)}},
$$

where $v=\left(v^{(1)}, v^{(2)}\right)$ and outside of the discontinuities of $y$. For $i=1, \ldots, N$ we have

$$
\begin{equation*}
\xi_{i}(t)=\bar{\xi}_{i} \text { and } l_{j} \cdot\left(\Delta_{i} v+\Delta_{i} y_{x} \xi_{i}\right)=0 j \neq k_{i} . \tag{2.10}
\end{equation*}
$$

along each line of discontinuity $x_{i}(t)$ where $y$ has a discontinuity in the $k_{i}$ th characteristic family. Here, $\Delta_{i} v=v\left(x_{i}(t)+, t\right)-v\left(x_{i}(t)-, t\right)$ and $l_{j}$ denotes the $j$ th left eigenvector of the matrix $\left(\begin{array}{cc}0 & 1 \\ a^{2} & 0\end{array}\right)$. The proof of Lemma 2.1 follows from Theorem $2.2[8]$ which has been repeated in the Appendix 5 for convenience. We observe that the system (2.2) fulfills the hypothesis (H1) to (H3) due to its linearity in the flux and since $f \in C^{4}$. The hypothesis are also detailed in the appendix 5 . Due to the linearity of the hyperbolic flux the eigenvectors and eigenvalues are independent of $y$. Therefore, the shock sensitivities $\xi_{i}$ are constant in time leading to equation (2.10). The evolution equation for $v$
strongly simplifies, too, leading to equation (2.9). According to [12, Definition 2] $\bar{u} \in \mathcal{U}$ and $\bar{y}=(\bar{u}, f(\bar{u}))$ are piecewise Lipschitz with simple discontinuities. The equations (2.9) and (2.10) are particularly simple due to the linear transport in the hyperbolic relaxation. In particular, the equation for $\xi_{i}$ is solved without any effort, since $\xi_{i}$ is constant in time. Further, if we diagonalize (2.2) first and then apply the tangent vector calculus, the second equation (2.10) simplifies. We come back to this point later in the numerical scheme. It is important to note that in the previous result it is assumed that all variations $u_{\delta}$ (resp. $y_{\delta}$ ) posses the same number of discontinuities which in the case of problem (1.1) is unknown a priori. An interpretation of the second condition of equation (2.10) may be as follows: given a discontinuity in family $i$ the condition prevents a discontinuity at the same position in another family $j \neq i$.

Tangent vectors and their property (2.8) can be used to compute the variations of the cost functional (2.5).
Lemma 2.3. Assume the assertions of Lemma 2.1 hold true. Assume that $t_{0}>T$. Let $J$ be given by (2.5) and assume $I$ is sufficiently large. Then, the variation of $J$ with respect to a tangent vector $(v, \xi)$ for initial data $y(0, x)=\left(u_{0}(x), f\left(u_{0}(x)\right)\right.$ is given by

$$
\begin{align*}
\Delta_{(v, \xi)} J\left(y, y_{d}\right) & =2 \int \chi_{I}(x)\left(y^{(1)}(T, x)-y_{d}(x)\right) v^{(1)}(T, x) d x+ \\
& \sum_{i=1}^{N\left(u_{0}\right)}\left(\left(y^{(1)}\left(T, x_{i}+\right)-y_{d}\left(x_{i}+\right)\right)+\left(y^{(1)}\left(T, x_{i}-\right)-y_{d}\left(x_{i}-\right)\right)\right) \Delta_{i} y^{(1)}(T, \cdot) \xi_{i}(T) . \tag{2.11}
\end{align*}
$$

In other words, calling $\gamma(\delta)$ a curve generating the tangent vector $(v, \xi)$ and $y_{\delta}$ the solution for initial datum $\gamma(\delta)$, we have:

$$
\begin{equation*}
J\left(y_{\delta}, y_{d}\right)=J\left(y, y_{d}\right)+\delta \Delta_{(v, \xi)} J\left(y, y_{d}\right)+o(\delta) \tag{2.12}
\end{equation*}
$$

Note that the summation is on the discontinuities of $y^{(1)}$ at terminal time $T$ related by Lemma 2.1 to the number of discontinuities of $u_{0}$ (resp. $y_{0}$ ) reflecting the original problem. The proof is similar to [16, Proposition 1] and omitted.

Lemma 2.3 and equation (2.7) already suggest a numerical method for solving (1.1) with cost functional (2.5). Given some control $u_{0}$ and a stepsize $\rho>0$ we obtain a new control $\tilde{u}_{0}$ corresponding to smaller value of the cost functional $J$ by the following variation

$$
\begin{equation*}
\tilde{u}_{0}(x)=u_{0}(x)-\left(\rho v(0, x)-\sum_{i: \xi_{i}(0)>0}^{N\left(u_{0}\right)} \Delta_{i} u_{0} \chi_{\left[x_{i}, x_{i}+\rho \xi_{i}(0)\right]}+\sum_{i: \xi_{i}(0)<0}^{N\left(u_{0}\right)} \Delta_{i} u_{0} \chi_{\left[x_{i}+\rho \xi_{i}(0), x_{i}\right]}\right) \tag{2.13}
\end{equation*}
$$

and $\tilde{y}_{0}=\left(\tilde{u}_{0}, f\left(\tilde{u}_{0}\right)\right)$. Further, $v(0, x)$ is the solution at time $t=0$ to (2.9) for terminal data

$$
\begin{equation*}
v^{(1)}(T, x)=\left(y_{d}(x)-y^{(1)}(T, x)\right), v^{(2)}(T, x)=0 \tag{2.14}
\end{equation*}
$$

and $\xi_{i}(0)$ the solution to (2.10) with terminal data

$$
\begin{equation*}
\xi_{i}(T)=\left(\left(y_{d}\left(x_{i}+\right)-y^{(1)}\left(T, x_{i}+\right)\right)+\left(y_{d}\left(x_{i}-\right)-y^{(1)}\left(T, x_{i}-\right)\right)\right) \Delta_{i} y^{(1)}(T) \tag{2.15}
\end{equation*}
$$

Now, the given choice for $(v(T, x), \xi(T))$ yields:

$$
\begin{array}{r}
\Delta_{(v, \xi)} J\left(y, y_{d}\right)=-2 \int \chi_{I}(x)\left(y^{(1)}(T, x)-y_{d}(x)\right)^{2} d x \\
-\sum_{i=1}^{N\left(u_{0}\right)}\left(\left(y^{(1)}\left(T, x_{i}+\right)-y_{d}\left(x_{i}+\right)\right)+\left(y^{(1)}\left(T, x_{i}-\right)-y_{d}\left(x_{i}-\right)\right)\right)^{2}\left(\Delta_{i} y^{(1)}(T, \cdot)\right)^{2}
\end{array}
$$

Denote by $\tilde{y}$ the solution to equation (2.2) with initial data $\tilde{y}_{0}=\left(\tilde{u}_{0}, f\left(\tilde{u}_{0}\right)\right)$. The solution to initial data $\left(u_{0}, f\left(u_{0}\right)\right)$ is denoted by $y$. Then, according to (2.12) we obtain

$$
\begin{equation*}
J\left(\tilde{y}, y_{d}\right)<J\left(y, y_{d}\right) \tag{2.16}
\end{equation*}
$$

for $\rho$ sufficiently small.
Note that this requires to solve (2.9) backwards in time and to fulfill (2.10). The previous computations motivate a numerical scheme for approximately solving (1.1), (2.2) and (2.5).

Before stating the full discrete algorithm we reformulate and comment on some parts of the method. The system (2.2) is diagonalisable with eigenvalues $\lambda_{1,2}= \pm a$ and characteristic variables

$$
\begin{equation*}
\eta^{(1)}=y^{(2)}+a y^{(1)} \text { and } \eta^{(2)}=y^{(2)}-a y^{(1)} . \tag{2.17}
\end{equation*}
$$

Also, in view of condition (2.10) it is numerically advantageous to consider the minimization problem for $J$ in characteristic variables $\eta=\left(\eta^{(1)}, \eta^{(2)}\right)$. Furthermore, in view of (2.13) equation (2.9) will be solved backwards in time for given terminal data $v(T, x)$. We obtain for $\tilde{v}(t, x)=v(T-t, x)$ the system

$$
\tilde{v}(0, \cdot)=v(T, \cdot), \tilde{v}_{t}-\left(\begin{array}{cc}
0 & 1  \tag{2.18}\\
a^{2} & 0
\end{array}\right) \tilde{v}_{x}=\frac{1}{\epsilon}\binom{0}{\tilde{v}^{(2)}-f^{\prime}\left(y^{(1)}(T-t, \cdot)\right) \tilde{v}^{(1)}} .
$$

and eigenvalues $\lambda_{1,2}=\mp a$ and corresponding characteristic variables

$$
\varphi^{(1)}=\tilde{v}^{(2)}+a \tilde{v}^{(1)} \text { and } \varphi^{(2)}=\tilde{v}^{(2)}-a \tilde{v}^{(1)} .
$$

Let $I=[0,1], T>0, y_{d}$ and $\epsilon>0$ be given. The cost functional $J$ is given by equation (2.5). We discuss the numerical discretization of problem (2.19) using a first-order finite volume scheme with periodic boundary conditions. Note that we leave (for the moment) both components of the initial data $y_{0}$ subject to optimization.

$$
\begin{equation*}
\min _{y_{0}=\left(y_{0}^{(1)}, y_{0}^{(2)}\right)} J \text { sbj to }(2.2), y(0, x)=y_{0}(x), y(t, 1)=y(t, 0), x \in I, t \geq 0 . \tag{2.19}
\end{equation*}
$$

Fix $a^{2}$ such that the subcharacteristic condition (2.4) is fulfilled. Introduce an equidistant spatial grid $\left\{x_{i}\right\}_{i=0}^{N_{x}}$ on $I$ with $\Delta x=x_{i+1}-x_{i}$. We choose $\Delta t$ such that the CFL condition holds, i.e., $\Delta t|a|=\Delta x$, and denote by $t^{n}=\Delta t n$ for $n=0, \ldots, N_{t}$. We write $x_{i+\frac{1}{2}}=x_{i}+\frac{\Delta x}{2}$ and for simplicity assume $x_{N_{x}}=1$ and $t^{N_{t}}=T$. Also for notational convenience we denote by $x_{-\frac{1}{2}}=x_{N_{x}-\frac{1}{2}}$ and $x_{N_{x}+\frac{1}{2}}=x_{\frac{1}{2}}$. Let $\mathcal{T}^{-1} \in \mathbb{R}^{2 \times 2}$ be the transformation to characteristic variables (2.17), i.e.,

$$
\mathcal{T}^{-1}\left(\begin{array}{cc}
0 & 1 \\
a^{2} & 0
\end{array}\right) \mathcal{T}=\left(\begin{array}{cc}
a & 0 \\
0 & -a
\end{array}\right) \text { and } \eta=\mathcal{T}^{-1} y
$$

On a time-interval $\left[t^{n}, t^{n+1}\right]$ we may introduce an operator splitting [25] to discretize source and transport term. In characteristic variables the splitting for $t \in\left[t^{n}, t^{n+1}\right]$ reads

$$
\begin{aligned}
& \partial_{t} \eta^{(1)}+a \partial_{x} \eta^{(1)}=0, \partial_{t} \eta^{(2)}=0, \\
& \partial_{t} \eta^{(1)}=S\left(\eta^{(1)}, \eta^{(2)}\right), \partial_{t} \eta^{(2)}=S\left(\eta^{(1)}, \eta^{(2)}\right), \\
& \partial_{t} \eta^{(2)}-a \partial_{x} \eta^{(2)}=0, \partial_{t} \eta^{(1)}=0 .
\end{aligned}
$$

where $S\left(\eta^{(1)}, \eta^{(2)}\right)=\frac{1}{\epsilon}\left(f\left((\mathcal{T} \eta)^{(1)}\right)-(\mathcal{T} \eta)^{(2)}\right)$. A discontinuity at time $t^{n}$ in any component of $\eta_{0}$ therefore moves with speed $a$ and $-a$, respectively. Clearly, we can also express the objective function $J$ in characteristic variables. The corresponding equations are given in the appendix 6 .

The cell average on $\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]$ at time $t^{n}$ for any function $u(t, x)$ is denoted by $u_{i}^{n}=\frac{1}{\Delta x} \int_{x_{i-\frac{1}{2}}}^{x_{i+\frac{1}{2}}} u\left(t^{n}, x\right) d x$. A first-order Upwind discretization of equation (2.2) using the reformulation in characteristic variables and an exact integration of the source term is given by (2.20) for $i=0, \ldots, N_{x}$ and $n=1, \ldots, N_{t}$.

$$
\begin{align*}
y_{i}^{0} & =\left(y_{0}\right)_{i}  \tag{2.20a}\\
\eta_{i}^{(1)} & =\left(\mathcal{T}^{-1} y_{i-1}^{n-1}\right)^{(1)}, \eta_{0}^{(1)}=\left(\mathcal{T}^{-1} y_{N_{x}}^{n-1}\right)^{(1)}, \eta_{i}^{(2)}=\left(\mathcal{T}^{-1} y_{i}^{n-1}\right)^{(2)},  \tag{2.20b}\\
\tilde{y}_{i}^{(1)} & =\left(\mathcal{T} \eta_{i}\right)^{(1)}, \tilde{y}_{i}^{(2)}=\exp \left(-\frac{\Delta t}{\epsilon}\right)\left(\mathcal{T} \eta_{i}\right)^{(2)}+\left(1-\exp \left(-\frac{\Delta t}{\epsilon}\right)\right) f\left(\tilde{y}_{i}^{(1)}\right),  \tag{2.20c}\\
\eta_{i}^{(2)} & =\left(\mathcal{T}^{-1} \tilde{y}_{i+1}\right)^{(2)}, \eta_{N_{x}}^{(2)}=\left(\mathcal{T}^{-1} \tilde{y}_{0}\right)^{(2)}, \eta_{i}^{(1)}=\left(\mathcal{T}^{-1} \tilde{y}_{i}\right)^{(1)},  \tag{2.20d}\\
y_{i}^{n} & =\mathcal{T} \eta_{i} . \tag{2.20e}
\end{align*}
$$

The discretization (2.20) uses a different splitting compared to [25] which leads to more complicated update formulas above but will be advantageous later on. In the current splitting we first transport in the first characteristic variable, then apply the source term and finally transport the second characteristic variable. The transformation to characteristic variables $\eta$ and the CFL condition allows to resolve the transport exactly. Due to the particular CFL chosen, after one time step any discontinuity located at a cell interface reaches after $\Delta t$ again a cell interfaces $x_{i+\frac{1}{2}}$.

As explained above, instead of equation (2.9) we discretize equation (2.18) to solve for the variations $\tilde{v}$. Similarly to (2.20), we transform (2.18) to characteristic variables $\varphi$ and resolve the linear transport exactly. If the discretized initial data is denoted by $\tilde{v}_{i}^{0}$, then we obtain for $i=0, \ldots, N_{x}$ and $n=1, \ldots, N_{t}$

$$
\begin{align*}
\varphi_{i}^{(2)} & =\left(\mathcal{T}^{-1} \tilde{v}_{i-1}^{n-1}\right)^{(2)}, \varphi_{0}^{(2)}=\left(\mathcal{T}^{-1} \tilde{v}_{N_{x}}^{n-1}\right)^{(2)}, \varphi_{i}^{(1)}=\left(\mathcal{T}^{-1} \tilde{v}_{i}^{n-1}\right)^{(1)}  \tag{2.21a}\\
\bar{v}_{i}^{(1)} & =\left(\mathcal{T} \varphi_{i}\right)^{(1)}, \bar{v}_{i}^{(2)}=\exp \left(-\frac{\Delta t}{\epsilon}\right)\left(\mathcal{T} \varphi_{i}\right)^{(2)}+\left(1-\exp \left(-\frac{\Delta t}{\epsilon}\right) f^{\prime}\left(y_{i}^{(1), N_{t}-n}\right) \bar{v}_{i}^{(1)}\right.  \tag{2.21b}\\
\varphi_{i}^{(1)} & =\left(\mathcal{T}^{-1} \bar{v}_{i+1}\right)^{(1)}, \varphi_{N_{x}}^{(1)}=\left(\mathcal{T}^{-1} \bar{v}_{0}\right)^{(1)}, \varphi_{i}^{(2)}=\left(\mathcal{T}^{-1} \bar{v}_{i}\right)^{(2)},  \tag{2.21c}\\
\tilde{v}_{i}^{n} & =\mathcal{T} \varphi_{i} . \tag{2.21d}
\end{align*}
$$

Next we turn to the discretization of equation (2.10). Within a first-order finite volume scheme a piecewise constant approximation is used to reconstruct the solution, i.e.,

$$
\begin{equation*}
y(t, x) \approx \sum_{i=0}^{N_{x}} \chi_{\left[t^{n}, t^{n+1}\right] \times\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]}(t, x) y_{i}^{n}, \tag{2.22}
\end{equation*}
$$

and similarly for the initial data. Therefore, numerically solving problem (2.19) naturally leads to consider piecewise constant controls $y_{0} \in \mathcal{U}$ having discontinuities at possibly each cell boundary $x_{i+\frac{1}{2}}$. Therefore, a shift in the position of the discontinuity $\xi_{i}$ may occur at each boundary $x_{i+\frac{1}{2}}$. As long as the spatial resolution is not modified the number of discontinuities is however fixed being a crucial assumption in Theorem 2.1. Since the conservative $y$ and characteristic variables $\eta$ are equivalent upon the linear transformation $\mathcal{T}$ we optimize in (2.19) for $\eta_{0}=\mathcal{T}^{-1} y_{0}$ instead of $y_{0}$.

For fixed $N_{x}$ the set of all admissible controls $\mathcal{U}_{a d} \subset \mathcal{U}$ consists of all piecewise constant functions $\eta_{0}(x)$ given by

$$
\begin{equation*}
\eta_{0}^{(j)}(x)=\sum_{i=0}^{N_{x}} \chi_{\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]}(x) \eta_{0, i}^{(j)}, j=1,2, \tag{2.23}
\end{equation*}
$$

which additionally fulfill (2.24), i.e.,

$$
\begin{equation*}
\left(\eta_{0}^{(1)}\right)_{2 i}=\left(\eta_{0}^{(1)}\right)_{2 i+1},\left(\eta_{0}^{(2)}\right)_{2 i-1}=\left(\eta_{0}^{(2)}\right)_{2 i}, i=0, \ldots, \frac{N_{x}}{2} \tag{2.24}
\end{equation*}
$$

Note that for $N_{x}$ sufficiently large the condition (2.24) still allows to approximate any piecewise constant function. The condition (2.24) allows as admissible controls piecewise constant $\eta_{0}^{(i)} \in \mathcal{U}$ for $i=1,2$ having each only $\frac{N_{x}}{2}$ points of discontinuity which are of the following type: the first component $\eta_{0}^{(1)}$ may have a discontinuity only at $x_{i+\frac{1}{2}}$ for some odd value $i$ and the second component $\eta_{0}^{(2)}$ may only have a discontinuity at $x_{i+\frac{1}{2}}$ for some even value $i$. The construction (2.23) and (2.24) allows that for any given discretization cell $i$, either the first or second component has a discontinuity across the cell boundary $i+\frac{1}{2}$. The other component is constant across the cell boundary. This ensures that the second part of condition (2.10) is fulfilled. Further, the jump is parallel to the eigenvectors and does not split under advection. This structure is preserved in the splitting scheme below and therefore condition (2.10) is automatically fulfilled after transport and application of the source term.

When computing the tangent vector to $\eta_{0}$ we now obtain the $L^{1}$-variations $\varphi_{0}$ and the variation in the position of the discontinuities $\xi_{i}$. We denote by $\xi_{i}, i=0, \ldots, N_{x}$ the variation of the discontinuity at position $x_{i+\frac{1}{2}}$. Hence, $\xi_{i}$ for $i$ odd (even) is the variation of the discontinuity in the first (second) component of $\eta_{0}$. In original variables the equation for $\xi_{i}$ is given by equation (2.15). Transformation in characteristic variables lead to equation (6.2c) repeated here for convenience:

$$
\xi_{j}(T)=\frac{1}{2}\left(\left(\bar{\Delta}_{i(j)} \eta^{(1), N_{t}}-\bar{\Delta}_{i(j)} \eta^{(2), N_{t}}-\bar{\Delta}_{i(j)} y_{d}\right)\right)\left(\widehat{\Delta}_{i(j)} \eta^{(1), N_{t}}-\widehat{\Delta}_{i(j)} \eta^{(2), N_{t}}\right)
$$

where $\bar{\Delta} w_{k}=\frac{1}{2 a}\left(w_{k+1}+w_{k}\right)$ and $\widehat{\Delta} w_{k}=w_{k+1}-w_{k}$. Here, $i(j)$ denotes the index of the location $x_{i+\frac{1}{2}}$ of the $j$ th discontinuity in the initial data, i.e., for given $j$ we determine $i$ such that $x_{i+\frac{1}{2}}=x_{j+\frac{1}{2}}+a T$. The equations for $\varphi_{0}$ are obtained from the equations (2.18) for $\tilde{v}$ and are given in detail in equation (6.2).

For the evolution of $\xi_{i}$ it is important that no new discontinuities are generated during the computation of transport and source term. This is guaranteed within each timestep $\Delta t$ using the scheme (2.20). To be more precise, the splitting does not introduce additional discontinuities due the action of the source term $S\left(\eta^{(1)}, \eta^{(2)}\right)$ provided that $\eta_{0}$ fulfills (2.24). Note that the same holds true for $\varphi$ provided that $\varphi_{i}^{0}=\mathcal{T}^{-1} \tilde{v}_{i}^{0}$ fulfills (2.24).

The evaluation of $\xi_{j}$ requires to compute the actual shock positions. Under assumption (2.24) the position of discontinuities in the first component and second component of $\eta$ at time $t^{n}$ are given by

$$
\begin{equation*}
x_{2 i-1}\left(t^{n}\right)=x_{2 i-1}(0)+a t^{n}, x_{2 i}\left(t^{n}\right)=x_{2 i}(0)-a t^{n}, i=0, \ldots, N_{x} . \tag{2.25}
\end{equation*}
$$

Here, $x_{j}(0)=x_{j+\frac{1}{2}}$ and we assume that discontinuities exiting at $x=1$ (or $x=0$ ) enter again at $x=0(x=1)$.

We turn to the discussion of the condition (2.10). In characteristic variables $\ell_{j}$ is the $j$ th unit vector. In the following we show that the second part of equation (2.10) is always fulfilled. Let $i$ be odd and $\eta$ be computed by the previous splitting. Then, we observe $\Delta_{i} \eta_{x}^{(2)}=0$ since $\eta^{(2)}$ is constant across the position of the discontinuity in the first family $x_{2 i-1}\left(t^{n}\right)$. Provided we discretize $J$ such that $\varphi_{j}^{0}:=\mathcal{T}^{-1} \tilde{v}_{j}^{0}, j=0, \ldots, N_{x}$ in (2.21) also fulfills (2.24), then, we have $\Delta_{i} \varphi^{(2)}=0$. The same is true for $i$ even and, therefore, the second part of equation (2.10) is fulfilled trivially within the proposed scheme. No additional modification of the scheme is required.

Finally, we consider $J\left(\eta, y_{d}\right)=J\left(\mathcal{T} y, y_{d}\right)$ and obtain the gradient of $J$ in terms of the characteristic variables $\eta$ and its associated tangent vectors $(\phi, \xi)$. The detailed equation for $\Delta_{v, \xi} J\left(\eta, y_{d}\right)$ is given in (6.1) with associated discretization given by equation (6.2).

In the following we describe an algorithm using tangent vectors to solve the optimization problem (2.19). In our description we may switch between characteristic $(\eta, \phi)$ and original variables $(y, v)$ by the linear and time-independent transformation $\mathcal{T}$ defined in equation (2.17). Further, $\tilde{v}$ and $v$ are related by $\tilde{v}(t, x)=v(T-t, x)$ and $\phi$ and $\varphi$ are also related by $\varphi(t, x)=\phi(T-t, x)$. We solve problem (2.19) by an iterative scheme starting with an initial guess for the control $u_{0} \in U_{a d}$. Herein, $U_{a d}$ is such that $\mathcal{T}^{-1}\left(u_{0}, f\left(u_{0}\right)\right)=\eta_{0}$ fulfills equation (2.24). Due to the discretization we have $N_{x}$ discontinuities in $\eta_{0}$ located at the cell interfaces $x_{i+\frac{1}{2}}$ for all $i$. With each discontinuity we have an associated shift $\xi_{i}$ (independent of time due to equation (2.10)). Further, we have a (discrete) solution $y$ (resp. $\eta$ ) associated with the control $\eta_{0} . \eta_{0}$ might not be optimal and should improved. The key to improve the control is equation (2.16) in original variables (see equation (6.1) for a formulation of the relevant term in characteristic variables). In order to fulfill the descent condition (2.16) we choose ( $v, \xi$ ) (resp. $(\varphi, \xi))$ such that

$$
\begin{equation*}
\Delta_{(v, \xi)} J\left(y, y_{d}\right) \equiv \Delta_{(\varphi, \xi)} J\left(y, y_{d}\right)<0 . \tag{2.26}
\end{equation*}
$$

This choice in original variables is given by equation (2.14) and (2.15). In characteristic variables it is given by equation (6.2) and (6.2c). Note that in the formulation of $\Delta_{(v, \xi)} J\left(y, y_{d}\right)$ we prescribe terminal data $t=T$ for the $L^{1}$-variation $v$ and the the shock shifts $\xi$. However, our control $u_{0}$, resp. $\eta_{0}$, is prescribed at initial time $t=0$. Therefore, we need to transport the particular choice of $v(T)$, resp. $\phi(T)$, such that (2.26) holds backwards in time using the (discrete) dynamics (2.18), resp. (2.21). For simplicity the backwards in time dynamics is transformed into a forward in time dynamics for the variables $\tilde{v}$, resp. $\varphi$. In case of $\xi_{i}(T)$ the transport backwards in time is trivial, since $\xi_{i}$ is constant in time. Now, we obtain a new control $u_{0}$, resp. $\eta_{0}$ by applying formula (2.13). In the numerical scheme we update the characteristic variable $\eta_{0}(x)$ by updating the cell average $\eta_{0, i}$ of each component at every spatial point $i=0, \ldots, N_{x}$. We proceed in two steps. First, we add componentwise the cell average of the $L^{1}$-part $\varphi_{i}^{N_{t}}=\phi_{i}^{0}$ in each cell $i$. Second, we need to take into account the variation $\xi_{j}(0)$ for each discontinuity $j$. In characteristic variables the location of the discontinuity $j$ is computed according
to equation (2.25) and it corresponds to the cell index $i(j)$. According to equation (2.23) the solution is piecewise constant on each cell $\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]$. Further, it has only discontinuities at every second cell interfaces in each component. In order to avoid interacting shock shifts we restrict $\rho$ in equation (2.13) such that $\left|\xi_{i}\right| \leq 2 \Delta x$. Then, we reconstruct $\eta_{0}$ according to (2.23) and apply the shifts by adding

$$
\sum_{i: \xi_{i}(0)>0}^{N\left(u_{0}\right)} \Delta_{i} u_{0} \chi_{\left[x_{i}(0), x_{i}(0)+\rho \xi_{i}(0)\right]}+\sum_{i: \xi_{i}(0)<0}^{N\left(u_{0}\right)} \Delta_{i} u_{0} \chi_{\left[x_{i}(0)+\rho \xi_{i}(0), x_{i}(0)\right]}
$$

as in equation (2.13). In fact, we add the corresponding term in characteristic variables. We further have to ensure that the new obtained $\eta_{0}$ still fulfills (2.24). Therefore, we need to compute cell averages on $\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]$ of the updated control $\eta_{0}$ such that (2.24) holds. Since $\eta_{0}$ was piecewise linear before the new cell averages are computed exactly.

In the following we present the update formula in characteristic variables. Denote the volume averaged shifted control in cell $j$ by $\Xi_{j}^{k} \Delta x$ for the $k$ th component. Then, we obtain

$$
\begin{align*}
& \Xi_{j-1}^{1}= \Xi_{j}^{1}=  \tag{2.27a}\\
& \min \left\{\left(-\bar{\xi}_{j}\right)^{+}, \Delta x\right\} \eta_{j+1}^{(1), 0}+\max \left\{\bar{\xi}_{j-2}^{+}-\Delta x, 0\right\} \eta_{j-1}^{(1), 0}+ \\
&\left(\Delta x-\min \left\{\left(-\bar{\xi}_{j}\right)^{+}, \Delta x\right\}-\max \left\{\bar{\xi}_{j-2}^{+}-\Delta x, 0\right\}\right) \eta_{j}^{(1), 0},  \tag{2.27b}\\
& \Xi_{k-1}^{2}= \Xi_{k}^{2}=\min \left\{\left(-\bar{\xi}_{k}\right)^{+}, \Delta x\right\} \eta_{k+1}^{(2), 0}+\max \left\{\bar{\xi}_{k-2}^{+}-\Delta x, 0\right\} \eta_{k-1}^{(2), 0}+ \\
& \quad\left(\Delta x-\min \left\{\left(-\bar{\xi}_{k}\right)^{+}, \Delta x\right\}-\max \left\{\bar{\xi}_{k-2}^{+}-\Delta x, 0\right\}\right) \eta_{k}^{(2), 0},  \tag{2.27c}\\
& \tilde{\eta}_{i}^{(1), 0}= \frac{\Xi_{i}^{1}}{\Delta x}-\varphi_{i}^{(1), N_{t}},  \tag{2.27d}\\
& \tilde{\eta}_{i}^{(2), 0}= \frac{\Xi_{i}^{2}}{\Delta x}-\varphi_{i}^{(2), N_{t}} .
\end{align*}
$$

Herein, the current cell average of the control is denoted by $\eta_{i}^{0}$ and the update is denoted by $\tilde{\eta}_{i}^{0}$. The index $j \in\left\{0, \ldots, N_{x}\right\}$ is odd, and $k \in\left\{0, \ldots, N_{x}\right\}$ is even. Further, we denote by $x^{+}=\max \{x, 0\}$ and by $\bar{\xi}_{i}=\mathcal{P}\left(-\xi_{i}\right)$ where $\mathcal{P}$ is the projection on $[-2 \Delta x, 2 \Delta x]$, i.e.,

$$
\mathcal{P}(z)=\left(\begin{array}{cc}
-2 \Delta x & z \leq-2 \Delta x \\
z & -2 \Delta x<z<2 \Delta x \\
2 \Delta x & z \geq 2 \Delta x
\end{array}\right)
$$

Those computations are exemplified in Remark 2.2. Summarizing, in (2.27) a piecewise constant reconstruction of $\tilde{\eta}$ is computed where for example in the case of the first component the discontinuity at $x_{i+\frac{1}{2}}$ is moved by $\xi_{i}$ and at $x_{i-\frac{3}{2}}$ by $\xi_{i-2}$. Since $\varphi$ fulfills (2.24) this holds true for $\tilde{\eta}$.

The previous computation leads to an iterative algorithm for numerically solving (2.19) or equivalently

$$
\min _{\eta_{0} \in \mathcal{U}} J\left(\mathcal{T} \eta, y_{d}\right) \text { sbj to } y=\mathcal{T} \eta,(2.2), \eta(0, x)=\eta_{0}(x), \eta(t, 0)=\eta(t, 1) \text { and }(2.24)
$$

## Algorithm

1. Set terminal time $T>0, a^{2} \geq \max _{y^{(1)}}\left(f^{\prime}\left(y^{(1)}\right)^{2}\right.$ and choose an equidistant spatial discretization with $N_{x}$ gridpoints. Choose $\Delta t=\frac{\Delta x}{a}$ and $\mathbf{k}=0$. Let $\eta_{0, i}^{\mathbf{k}}=$ $\left(\left(\eta_{0, i}^{(1)}\right)^{\mathbf{k}},\left(\eta_{0, i}^{(2)}\right)^{\mathbf{k}}\right)$ for $i=0, \ldots, N_{x}$, be an arbitrary initial control such that $\eta_{0}^{\mathbf{k}}$ fulfills (2.24). Let $\left(y_{d}\right)_{i}$ be a discretization of the given function $y_{d}(\cdot)$.
2. Solve equations (2.20) with initial data $\left(y_{0}\right)_{i}:=\mathcal{T} \eta_{0, i}^{\mathbf{k}}$ to obtain $\eta_{i}^{N_{t}}=\mathcal{T}^{-1} y_{i}^{N_{t}}$.
3. Set initial data $\varphi_{i}^{0}$ and shock variations $\xi_{i}$ given by equation (6.2). Therein, $x_{i}\left(t^{N_{t}}\right)$ are given by equation (2.25).
4. Solve equations (2.21) for initial data $\tilde{v}_{i}^{0}:=\mathcal{T} \varphi_{i}^{0}$ to obtain $\varphi_{i}^{N_{t}}=\mathcal{T}^{-1} \tilde{v}_{i}^{N_{t}}$.
5. Obtain the new iterate $\eta_{0, i}^{\mathbf{k}+\mathbf{1}}:=\tilde{\eta}_{i}^{0}$ where $\tilde{\eta}_{i}^{0}$ is given by equation (2.27). To evaluate equation (2.27) use the old iterate as $\eta_{i}^{0}:=\eta_{0, i}^{\mathbf{k}}$, the solution $\varphi_{i}^{N_{t}}$ and the shifts $\xi_{i}$.
6. Provided that $J\left(\mathcal{T} \eta, y_{d}\right)$ is sufficiently small we terminate. Otherwise set $\mathbf{k} \rightarrow$ $\mathbf{k}+1$ and continue with step (2).
Remark 2.2. Note that $\eta_{0, i}^{\mathbf{k}}$ is a vector of two components. The algorithm computes a sequence $\left(\left(\eta_{0, i}^{\mathbf{k}}\right)_{i=0}^{N_{x}}\right)_{\mathbf{k}=0}^{\infty}$ of approximations

$$
y^{\mathbf{k}}(0, x)=\sum_{i=0}^{N_{x}} \chi_{\left[x_{i-\frac{1}{2}, x_{i+\frac{1}{2}}}\right]} \mathcal{T} \eta_{0, i}^{\mathbf{k}}
$$

to the optimal control $y^{*}(0, x)$ of problem (2.19). During the iteration on $\mathbf{k}$ the data $\varphi_{i}^{0}$ and $\xi_{i}$ are chosen in every step such that the gradient of the objective function is non-positive, see equation (6.1) and equation (6.2). The derivative of the objective functional is computed at terminal time $T$ using the solution $\eta$ at time $T$. The update for the control $\eta_{0, i}$ has to be at initial time. Therefore, in step 4, the algorithm propagates the information backwards using the equations (2.21). At initial time the variations in $L^{1}$ and for the shock position are added towards the old control $\eta_{0, i}^{\mathbf{k}}$ in order to obtain an improved control for the next iteration $\mathbf{k}+\mathbf{1}$. Using this new control the equation (2.20) will be solved. If the mismatch in the objective functional is still large enough the procedure is repeated.

We present more details on the motivation of equation (2.27) where the values $\Xi_{j}$ are introduced. Those values describe the effect of the shift in the discontinuity on the values of $\eta_{i}^{0}$ and are used in step 5 of the algorithm. We consider the first component only with the other being similar. Due to the construction of the discretized control $\eta_{i}^{(1), 0}$ we have a discontinuity only at indices $i$ when $i$ is odd. More precisely, the discontinuity is located at position $x_{i+\frac{1}{2}}$ for $i$ odd. From step 4 we obtain the shift $\xi_{i}$ in the shock position. We revert the sign of $\xi_{i}$ in order to obtain a descent in the objective function
(see equation for $\bar{\xi}_{i}$ below). Since we require equation (2.24) to hold we first project the shift $\xi_{i}$ onto the interval $[-2 \Delta x, 2 \Delta x]$ and denote the projection by $\bar{\xi}_{i}=\mathcal{P}\left(-\xi_{i}\right)$. Assume now $-\Delta x<\bar{\xi}_{i}<0$. Then, in order to have a piecewise constant control on the spatial intervals $\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]$ we need to shift the current control at point $x_{i+\frac{1}{2}}$ by $\xi_{i}$ and average again the shifted control on $\left[x_{i-\frac{1}{2}}, x_{i+\frac{1}{2}}\right]$. This way we obtain $\eta_{0, i}^{(1)}$. Assume no other shift has happened. Then, the volume average of the shifted control in cell $i$ is given by

$$
\eta_{0, i}^{(1)}=\frac{1}{\Delta x}\left(\left(-\bar{\xi}_{i}\right) \eta_{0, i+1}^{(1)}+\left(\Delta x-(-\bar{\xi}) \eta_{0 . i}^{(1)}\right)\right)=\frac{\Xi_{i}^{1}}{\Delta x}
$$

The formula equation (2.27) is obtained when taking into account shifts in positive and negative direction as well as shifts of the neighboring cells.

Note that in the case $\xi_{i} \equiv 0$ for all $i$, we obtain $\frac{\Xi_{i}^{1}}{\Delta x}=\eta_{i}^{(1), 0}$
Compared with the general case given by equation (2.13) we presented the algorithm for $\rho \equiv 1$. In the numerical results we observed convergence even in this case. However, it is only minor modification to include $0<\rho<1$ in step 5. Note that including a $\rho>0$ requires also to introduce a stepsize control mechanism as for example Armijo's rule.
3. Numerical results We present numerical results using the previous calculus for two cases. The simplest possible application is the optimal control of a linear system (2.2) without source term. Second, we present results on the optimal control for the relaxation system (2.2) for Burgers' flux $f(w)=\frac{1}{2} w^{2}$. All spatial grids are equidistant on $I=[0,1]$ and the temporal discretization is such that the CFL condition [15] is satisfied. We use periodic boundary conditions in all cases. The cost functional $J$ is discretized using the trapezoidal rule. All initial controls are constant with $\eta_{0}(x)=\left(\frac{1}{2}, \frac{1}{4}\right)$. Numerical tests not reported here show first-order convergence of the scheme. We observe that the applied modifications to the IMEX scheme do not alter the properties of the original scheme proposed in [25].

The minimization problem (1.1) does not have necessarily a unique solution. This is for example the case if the desired state is generated by a strong discontinuity. Therefore, the initial value for the optimization might have a strong influence on the local minimizer found. The non uniqueness of the minimization problem has been discussed in recent literature as for example [1]. Our purpose of studying problem (1.1) is to exemplify the numerical application of the tangent vector calculus. Therefore, we do not focus on possible non uniqueness of the solution to (1.1).
3.1. Optimal control of a linear system We consider problem (3.1) with periodic boundary conditions on the domain $x \in I, T=0.35$ and $a^{2}=\frac{5}{4}$.

$$
\min _{u_{0}} \int_{0}^{1}\left(y^{(1)}(T, x)-y_{d}^{(1)}(x)\right) d x \text { sbj to } y_{t}+\left(\begin{array}{rr}
0 & 1  \tag{3.1}\\
a^{2} & 0
\end{array}\right) y_{x}=0, y(0, x)=\left(u_{0}(x), f\left(u_{0}\right)\right) .
$$

In the example we initialize the control similar to (2.3) and $f(w)=\frac{1}{2} w^{2}$. The desired state $y_{d}$ is depicted in Figure 3.1 (along with the obtained optimized state). The pointwise error in terminal state are depicted in Figure 3.2. The corresponding initial data $y(0, x)=\left(u_{0}(x), f\left(u_{0}\right)\right)$ corresponding to the desired state $y_{d}$ is given by $u_{0}(x)=\left(\begin{array}{cc}1 & 0.30 \leq x \leq 0.45 \\ 0.1 & 0.5 \leq x \leq 0.65 \\ 0 & \text { else }\end{array}\right)$.


Fig. 3.1. Desired state (red cross) and optimized state (blue circles) for the linear system with $N_{x}=200$ grid points in space for both components $\left(y^{(1)}, y^{(2)}\right)$ in the left part of the figure. The pointwise difference of optimized and desired state is given in Figure 3.2. Iteration history in $\log -$ scale for the cost. The spatial resolution of the scheme is $\Delta x=5 \times 10^{-3}$ and $\log _{10}(\Delta x) \approx-2.3$.


Fig. 3.2. Pointwise difference of desired state and optimized state in logarithmic scale for the linear system with $N_{x}=200$ grid points in space for both components $\left(y^{(1)}, y^{(2)}\right)$. The corresponding profiles are depicted in Figure 3.1.

The algorithm is started with constant control $\eta_{0}$ and stopped after at most 2000 steps of iteration. The iteration history is depicted in the right part of Figure 3.1. With the proposed method we observe convergence until the grid resolution is reached. The dependence on the spatial grid is given in Table 3.1. As expected we observe first-order convergence.

We compare the obtained results also with the case when the shock variations $\xi_{i}$ are not taken into account. Hence, we consider the same example as before and run the same algorithm but setting $\xi_{i} \equiv 0$ for $j=1,2$ in equation (2.27). The dependence on the spatial grid in this situation is depicted in Table 3.2. Comparing with Table 3.1 we observe a deterioration in the convergence rate from $\approx 1$ when including the shock variations to $\approx \frac{1}{2}$ for the cost functional and $\approx \frac{3}{4}$ for state when neglecting this contribution.

Note that the error in the second component is smaller compared with the scheme neglecting the shock sensitivities. In the first component and the functional we observe the opposite behavior. However the overall resolution is always below the numerical resolution of the grid and, therefore, a comparison of those values is difficult.

| $N_{x}$ | $J$ | Rate | $\left\\|y_{0}^{(1)}-y_{*}^{(1)}(t=0, \cdot)\right\\|_{2}$ | Rate | $\left\\|y_{0}^{(2)}-y_{*}^{(2)}(t=0, \cdot)\right\\|_{2}$ | Rate |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 50 | $4.6936 e-03$ | $(0.00)$ | $5.5364 e-03$ | $(0.00)$ | $6.9206 e-03$ | $(0.00)$ |  |
| 100 | $2.5878 e-03$ | $(0.91)$ | $2.8922 e-03$ | $(0.96)$ | $3.6152 e-03$ | $(0.96)$ |  |
| 200 | $1.3510 e-03$ | $(0.96)$ | $1.4739 e-03$ | $(0.98)$ | $1.8424 e-03$ | $(0.98)$ |  |
| 400 | $9.9964 e-04$ | $(0.68)$ | $8.2176 e-04$ | $(0.90)$ | $1.2031 e-03$ | $(0.77)$ |  |
| $-N_{x}$ |  |  |  |  |  |  | $\left\\|y_{d}^{(1)}-y_{*}^{(1)}(t=T, \cdot)\right\\|_{2}$ |
|  | Rate | $\left\\|y_{d}^{(2)}-y_{*}^{(2)}(t=T, \cdot)\right\\|_{2}$ | Rate |  |  |  |  |
|  | 50 | $5.5364 e-03$ | $(0.00)$ | $6.9206 e-03$ | $(0.00)$ |  |  |
|  | 100 | $2.8922 e-03$ | $(0.96)$ | $3.6152 e-03$ | $(0.96)$ |  |  |
|  | 200 | $1.4739 e-03$ | $(0.98)$ | $1.8424 e-03$ | $(0.98)$ |  |  |
|  | 400 | $7.5945 e-04$ | $(0.97)$ | $1.2666 e-03$ | $(0.73)$ |  |  |

Table 3.1. Convergence history for different spatial grids with $N_{x}$ equidistant grid points. Reported are the value of the cost functional J after optimization, the $L^{2}$-norm difference of both components $\left(y^{(1)}, y^{(2)}\right)$ for the initial data $y_{0}=y(0, x)$ and the desired state $y_{d}$. The optimized solution is denoted by $y_{*}(t, x)$.

| $N_{x}$ | $J$ | Rate | $\left\\|y_{0}^{(1)}-y_{*}^{(1)}(t=0, \cdot)\right\\|_{2}$ | Rate | $\left\\|y_{0}^{(2)}-y_{*}^{(2)}(t=0, \cdot)\right\\|_{2}$ | Rate |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 50 | $9.8550 e-04$ | $(0.00)$ | $1.9823 e-03$ | $(0.00)$ | $3.6662 e-03$ | $(0.00)$ |
| 100 | $9.8310 e-04$ | $(0.50)$ | $1.3402 e-03$ | $(0.74)$ | $2.6691 e-03$ | $(0.69)$ |
| 200 | $9.9839 e-04$ | $(0.49)$ | $9.5812 e-04$ | $(0.70)$ | $1.8837 e-03$ | $(0.71)$ |
| 400 | $9.9986 e-04$ | $(0.50)$ | $6.7520 e-04$ | $(0.71)$ | $1.3360 e-03$ | $(0.70)$ |
| $N_{x}$ | $\left\\|y_{d}^{(1)}-y_{*}^{(1)}(t=T, \cdot)\right\\|_{2}$ | Rate | $\left\\|y_{d}^{(2)}-y_{*}^{(*)}(t=T, \cdot)\right\\|_{2}$ | Rate |  |  |
| 50 | $1.6294 e-03$ | $(0.00)$ | $3.9955 e-03$ | $(0.00)$ |  |  |
| 100 | $1.0328 e-03$ | $(0.79)$ | $2.8747 e-03$ | $(0.69)$ |  |  |
| 200 | $7.3280 e-04$ | $(0.70)$ | $2.0440 e-03$ | $(0.70)$ |  |  |
| 400 | $4.9368 e-04$ | $(0.74)$ | $1.4578 e-03$ | $(0.70)$ |  |  |

Table 3.2. Convergence history for different spatial grids with $N_{x}$ equidistant grid points neglecting the effect of the shock variations $\xi_{i}$ on the optimal control. As in Table 3.1 we report the value of the cost functional $J$ after optimization, the $L^{2}$ - norm difference of both components $\left(y^{(1)}, y^{(2)}\right)$ for the initial data $y_{0}=y(0, x)$ and the desired state $y_{d}$. The optimized solution is denoted by $y_{*}(t, x)$.

### 3.2. Optimal control of the relaxation approximation to Burgers' Equation

We consider problem (3.1) with $T=0.35, \epsilon=10^{-4}, N_{x}=400$ and $a=\frac{5}{4}$. The desired state $y_{d}$ is the solution at time $T$ to (2.2) and (2.3) for $u_{0}(x)=\cos (2 \pi x)$. The algorithm is started with constant control $\eta_{0}$ and stopped after the value of the cost functional is below $10^{-4}$. We depict desired and optimized state in the original variables $y(T, \cdot)$ in Figure 3.3. The obtained optimized control $\eta_{0}^{\mathbf{k}}$ and the function $u_{0}$ are depicted in Figure 3.3 and the iteration history in Figure 3.4.



Fig. 3.3. Desired state (red cross) and optimized state (blue circles) for the linear system with $N_{x}=400$ grid points in space for the variable $y^{(1)}(T, x)$ in the left part of the figure and for $y^{(1)}(0, x)$ in the right part.


FIg. 3.4. Iteration history for the example of Figure 3.3 in $\log$-scale for the value of the cost functional $J$. The stopping criteria is $J \leq 10^{-4}$ which is comparable with the spatial resolution of the scheme, $\Delta x=2.5 \times 10^{-3}$.

We observe a good agreement in the recovered desired state $y_{d}$. The difference in the obtained control is due to the fact that the solution to equation (2.19) is not unique.

In Figure 3.5 and Figure 3.6 we present a similar example but with desired state $y_{d}$ obtained as solution to (2.2) and (2.3) for $u_{0}(x)=\chi_{[0.3,0.45]}(x)$. All other parameters are as before. Similarly to the previous results we observe that there is no uniqueness in the control $u_{0}$ leading to $y_{d}$.

## 4. Summary



Fig. 3.5. Desired state (red cross) and optimized state (blue circles) for the linear system with $N_{x}=400$ grid points in space for the variable $y^{(1)}(T, x)$ in the left part of the figure and for $y^{(1)}(0, x)$ in the right part.


Fig. 3.6. Iteration history for the example of Figure 3.3 in $\log$-scale for the value of the cost functional J. The stopping criteria is $J \leq 10^{-4}$ which is comparable with the spatial resolution of the scheme is $\Delta x=2.5 \times 10^{-3}$.

We present a numerical method for solving optimal control problems subject to the relaxation approximation to hyperbolic conservation laws. Equations for the evolution of the corresponding tangent vector of the relaxation system are derived and a numerical discretization has been introduced. The tangent vector has been used to compute the analytical gradient of the reduced cost also in the presence of traveling discontinuities. A numerical discretization of the gradient has been implemented to solve some examples of linear and nonlinear optimal control problems. The computation of the finite dimensional part of the tangent vector allows to obtain the expected order of convergence.

We comment on Table 3.1 and Table 3.2. Compared with the method ignoring the shifts in the tangent vectors we observed an improved numerical performance. However, the improvement yields only higher convergence rates, but not necessarily smaller residuals (up to the chosen tolerance). This is not contradictory to the theoretical result. Therein, we only to expect residuals as low as the numerical discretization error. The iterative procedure using shifts does not yield the same iterates and includes a projection towards the piecewise constant control. This introduces an error of the order of the grid which may explain the higher residuals. In the linear example where the exact solution
is known, this solution can be obtained using either the shifts or suitable updates of the cell averages and ignoring the shifts. This is due to the linear structure of the problem. But, neglecting shifts yields a different approximation to the optimization problem. In particular, this is not the theoretical correct one. Nevertheless, it yields smaller numerical residuals. A possible reason might be that when ignoring the shifts we update on $N_{x}$ cells cell averages. The true solution however requires only two correct shock positions. Therefore, there are more degrees of freedom leading possibly to smaller residuals.
Acknowledgments. This work has been supported by DFG Cluster of Excellence EXC128, RWTH Seedfund project, the BMBF KinOpt Project and NSF KI-Net.

## 5. Appendix on definitions and supplementary lemmas

In this section we collect definition and statements from the reference $[8,10]$. They are given for sake of completeness.

Definition 5.1 (Continuous path). A mapping $\gamma:[a, b] \rightarrow L^{1}\left(\mathbb{R}^{n}\right)$ is called a continuous path, if $\gamma$ is continuous on the interval $[a, b]$ with respect to $L^{1}$-norm, i.e.,

$$
\forall x \in[a, b]: \lim _{\epsilon \rightarrow 0}\|\gamma(x+\epsilon)-\gamma(x)\|_{L^{1}}=0
$$

Definition 5.2 (Broad solution). Consider the quasi-linear partial differential equation

$$
\begin{equation*}
u_{t}(t, x)+A(t, x) u_{x}(t, x)=h(t, x, u), \tag{5.1}
\end{equation*}
$$

where $A \in \mathbb{R}^{n \times n}$ is strictly hyperbolic, Lipschitz and $h$ is measurable w.r.t. ( $t, x$ ) and Lipschitz continuous w.r.t. u. Assume an initial condition $u(0, x)=u^{0}(x)$ with $u_{0} \in$ $L^{1}\left(\mathbb{R} ; \mathbb{R}^{n}\right)$. Denote by $\ell_{i}, r_{i}$ the ith left and right eigenvectors of $A$. Denote by $\lambda_{i}$ the $i$ ith eigenvalues of $A$. We denote by $t \rightarrow y_{i}(t ; \tau, \xi)$ the solution to the Cauchy problem

$$
\frac{d}{d t} y(t)=\lambda_{i}(t, y(t)), y(\tau)=\xi
$$

Denote by $<,>$ the scalar product on $\mathbb{R}^{n}$ and by

$$
g_{i}:=<\ell_{i}, h>+<\partial_{t} \ell_{i}+\lambda_{i} \partial_{x} \ell_{i}, u>, u=\sum u_{i} r_{i} .
$$

We define a broad solution $u=\sum u_{i} r_{i}$ to equation (5.1) as a locally integrable function fulfilling

$$
\frac{d}{d t} u_{i}\left(t, y_{i}(t ; \tau, \xi)\right)=g_{i}\left(t, y_{i}(t ; \tau, \xi), u\left(t, y_{i}(t ; \tau, \xi)\right)\right.
$$

in the sense that for a.e. $(\tau, \xi)$ and all $i=1, \ldots, n$ the following holds

$$
u_{i}(\tau, \xi)=u_{i}^{0}\left(y_{i}(0 ; \tau, \xi)\right)+\int_{0}^{\tau} g_{i}\left(s, y_{i}(s ; \tau, \xi), u\left(s, y_{i}(s ; \tau, \xi)\right) d s\right.
$$

The main result used in this work is [8, Theorem 2.2]. We recall the statement for convenience.

Consider the equation

$$
\begin{equation*}
\partial_{t} u+\partial_{x} F(u)=h(t, x, u) . \tag{5.2}
\end{equation*}
$$

supplemented with initial data $u(0, x)=u_{0}(x)$ and the assumptions
(H1) The vector field $F: \Omega \rightarrow \mathbb{R}^{n}$ is $\mathcal{C}^{2}$ where $\Omega \subset \mathbb{R}^{n}$ is closed and bounded. For each $u \in \Omega$ the matrix $A(u)=D F(u)$ has $n$ real distinct eigenvalues. Its eigenvalues $\lambda_{i}$ and its left and right eigenvectors $\ell_{i}$ and $r_{i}$, respectively, are normalized such that $\left\langle\ell_{i}, r_{j}\right\rangle=\delta_{i j}$. Denote by

$$
A(u, v)=\int_{0}^{1} A(\theta u+(1-\theta) v) d \theta
$$

with corresponding eigenvectors $\ell_{i}(u, v), r_{i}(u, v)$ and eigenvalues $\lambda_{i}(u, v)$. Suppose that $\ell_{i}(u, v), r_{i}(u, v)$ and $\lambda_{i}(u, v)$ are uniformly bounded for all $u, v \in \Omega$.
(H2) Denote by $\hat{\lambda}$ the uniform bound on $\lambda_{i}(i, v)$ for all $i$. Then, solutions to (5.2) are considered in the domain

$$
\mathcal{D}:=\{(t, x): 0 \leq t \leq T, x \in[a+\hat{\lambda} t, b-\hat{\lambda} t]\}
$$

Assume further that the function $h: \mathcal{D} \times \Omega \rightarrow \mathbb{R}^{n}$ is bounded and continuously differentiable.
(H3) Whenever $u^{+} \in \Omega$ and $u^{-} \in \Omega$ are connected by a shock or a contact discontinuity, say of the $k$ th characteristic family, the linear system

$$
\begin{array}{r}
0=\Phi_{i}\left(u^{+}, u^{-}, w^{+}, w^{-}\right)=\sum_{j=1}^{n}<D \ell_{i}\left(u^{+}, u^{-}\right) \cdot\left(w_{j}^{+} r_{j}^{+}, w_{j}^{-} r_{j}^{-}\right), u^{+}-u^{-}>+ \\
\\
\sum_{j=1}^{n}<\ell_{i}\left(u^{+}, u^{-}\right), w_{j}^{+} r_{j}^{+}-w_{j}^{-} r_{j}^{-}>, \forall i \neq k
\end{array}
$$

can be uniquely solved in terms of the outgoing variables $w_{j^{ \pm}}^{ \pm} j^{ \pm} \in\left\{j^{-}: j<\right.$ $k\} \cup\left\{j^{+}: j>k\right\}=: \mathcal{O}$. Assume that the function $W_{j}$ defined by

$$
w_{j}^{ \pm}=W_{j^{ \pm}}\left(u^{+}, u^{-}\right)\left(\left(w_{j}\right)_{j^{ \pm} \notin \mathcal{O}}\right), j \neq k, j^{ \pm} \in \mathcal{O}
$$

satisfies a bound of the form

$$
\left\|W_{j^{ \pm}}\left(u^{+}, u^{-}\right)\left(\left(w_{j}\right)_{j^{ \pm} \notin \mathcal{O}}\right)\right\| \leq C\left\|\left(w_{j}\right)_{j^{ \pm} \notin \mathcal{O}}\right\|
$$

Here, $r_{j}^{ \pm}=r_{j}\left(u^{ \pm}\right)$. For a definition of the class of functions which are piecewise Lipschitz with simple discontinuities we refer to [8].
Theorem 5.3. Let the assumptions (H1)-(H3) hold true. Let u be a piecewise Lipschitz continuous solution to equation (5.2) with $u^{0}$ in the class PLSD. Let $\left(v_{0}, \xi_{0}\right) \in L^{1} \times \mathbb{R}^{N}$ be a tangent vector to $u^{0}$ generated by a regular variation $\gamma: \delta \rightarrow u_{\delta}^{0}$, Let $u_{\delta}$ be the solution of equation (5.2) with initial condition $u_{\delta}^{0}$. Then, there exists $\tau_{0}>0$ such that for all $t \in\left[0, \tau_{0}\right]$ the path $\bar{\gamma}: \delta \rightarrow u_{\delta}$ is a regular variation for $u_{\delta}(t, \cdot)$ generating the tangent vector $(v(t), \xi(t)) \in L^{1} \times \mathbb{R}^{N}$. The vector is the unique broad solution of the initial boundary value problem

$$
\begin{array}{r}
\xi(0)=\xi_{0}, v(0, x)=v_{0}(x), \\
v_{t}+A(u) v_{x}+(D A(u) v) u_{x}=h_{u}(t, x, u) v,
\end{array}
$$

outside the discontinuities of $u$ while for $\alpha=1, \ldots, N$

$$
\begin{array}{r}
<D \ell_{i}\left(u^{+}, u^{-}\right) \cdot\left(v^{+}+\xi_{\alpha} u_{x}^{+}, v^{-}+\xi_{\alpha} u_{x}^{-}\right), u^{+}-u^{-}>= \\
+<\ell_{i}\left(u^{+}, u^{-}\right), v^{+}+\xi_{\alpha} u_{x}^{+}-v^{-}-\xi_{\alpha} u_{x}^{-}>, i \neq k_{\alpha} \\
\frac{d}{d t} \xi_{\alpha}=D \lambda_{k_{\alpha}}\left(u^{+}, u^{-}\right)\left(v^{+}+\xi_{\alpha} u_{x}^{+}, v^{-}+\xi_{\alpha} u_{x}^{-}\right)
\end{array}
$$

along each line $x=x_{\alpha}(t)$ where $u$ suffers a discontinuity in the $k_{\alpha}$ characteristic direction.

The technical details are given in [8]. We motivate the result by the following explanation. The equation for $v$ is independent of the equation for $\xi_{\alpha}$ and can therefore be solved separately. The equation is formally derived by differentiation of the conservation law (5.2) with respect to $u$ in direction $v$. It describes the sensitivity of the $L^{1}$-part of the solution. The equation is linear in $v$, but not in conservative form and with possibly discontinuous coefficient due to discontinuities in $u$. Along those discontinuities the equations for $\xi_{\alpha}$ and the algebraic condition needs to hold true.

Assume $u$ is discontinuous across $\left(t, x_{\alpha}(t)\right)$. Then, for $u$ to be a weak solution the Rankine-Hugoniot condition needs to be fulfilled. The condition is $s[u]=[F(u)]$. Denoting the left and right sided limits $u\left(t, x_{\alpha}(t) \pm\right)$ by $u \pm$, respectively, the previous equation may be rewritten as

$$
\lambda\left(u^{+}, u^{-}\right)\left(u^{+}-u^{-}\right)=A\left(u^{+}, u^{-}\right)\left(u^{+}-u^{-}\right)
$$

using the average matrix $A\left(u^{+}, u^{-}\right)$introduced in (H1) and using that the shock speed $s$ is equal to an eigenvalue $\lambda\left(u^{+}, u^{-}\right)$of the matrix $A\left(u^{+}, u^{-}\right)$. If $x_{\alpha}$ is a discontinuity in the $k_{\alpha}$ family then the previous equation can be reduced to

$$
\frac{d}{d t} x_{\alpha}(t)=\lambda_{k_{\alpha}}\left(u^{+}, u^{-}\right),<\ell_{i}\left(u^{+}, u^{-}\right), u^{+}-u^{-}>=0 \forall i \neq k_{\alpha}
$$

According to the notion of derivatives we need to consider a variation not only in the $L^{1}$ part, but also in the shock position. Formally, this corresponds to considering the $O(\epsilon)$ perturbation of $u^{ \pm}=u\left(t, x_{\alpha}(t)\right)$ as

$$
u_{\epsilon}^{ \pm}=u\left(t, x_{\alpha}(t) \pm \epsilon \xi_{\alpha}(t)\right)+\epsilon v\left(t, x_{\alpha}(t)\right) .
$$

Assuming the perturbation $x_{\alpha}(t)+\epsilon \xi_{\alpha}(t)$ also fulfills the previous set of equations we may compute the formal limit for $\epsilon \rightarrow 0$. Hence, we obtain up to $O\left(\epsilon^{2}\right)$

$$
\frac{d}{d t} \xi_{\alpha}(t)=D \lambda_{k_{\alpha}}\left(u^{+}, u^{-}\right)\left(u_{x}^{+} \xi_{\alpha}+v^{+}, u_{x}^{-} \xi_{\alpha}+v^{-}\right)
$$

and similarly for the second algebraic equation. The second algebraic condition is interpreted as follows: a shift variation in the $k_{\alpha}$ family leads to a modified function which only suffers discontinuities in the same $k_{\alpha}$ family, but not in any other family $i \neq k_{\alpha}$.

## 6. Appendix on the expression of objective and tangent vectors in characteristic variables

The gradient of the objective functional $J$ in terms of characteristic variables $\eta$ and the associated tangent vector $(\phi, \xi)$ is given by

$$
\begin{align*}
\Delta_{(\phi, \xi)} J\left(\eta, y_{d}\right) & =\frac{1}{2} \int \chi_{I}(x)\left(\frac{\eta^{(1)}(T, x)-\eta^{(2)}(T, x)}{2 a}-y_{d}(x)\right) \phi^{(1)}(T, x) d x+  \tag{6.1}\\
& -\frac{1}{2} \int \chi_{I}(x)\left(\frac{\eta^{(1)}(T, x)-\eta^{(2)}(T, x)}{2 a}-y_{d}(x)\right) \phi^{(2)}(T, x) d x+ \\
& \frac{1}{2} \sum_{j=1}^{N\left(u_{0}\right)}\left\{\left(\frac{\eta^{(1)}\left(T, x_{j}(T)+\right)-\eta^{(2)}\left(T, x_{j}(T)+\right)}{2 a}-y_{d}\left(x_{j}(T)+\right)\right)+\right. \\
& \left.\left(\frac{\eta^{(1)}\left(T, x_{j}(T)-\right)-\eta^{(2)}\left(T, x_{j}(T)-\right)}{2 a}-y_{d}\left(x_{j}(T)-\right)\right)\right\} \times \xi_{j}(T) \times \\
& \left(\Delta_{j} \eta^{(1)}(T, \cdot)-\Delta_{j} \eta^{(2)}(T, \cdot)\right)
\end{align*}
$$

where $\Delta_{j} w(\cdot)=w\left(x_{j}(T)+\right)-w\left(x_{j}(T)-\right)$. Note that in the last terms the evaluation of $\eta$ at time $T$ is at point $x_{j}(T)$. This point is the terminal point of the $j$ th discontinuity in the initial data $\eta_{0}$ propagated forward to terminal time $T$. The value of $x_{j}$ is computed by moving the $j$ th discontinuity with speed $\pm a$ depending on whether it is a jump in the first or second component, see equation (2.25). For example, in equation (6.1) we have $x_{j}=x_{j}(0)-a T$ for all even indices $j$.

Since $\varphi(t, x)=\phi(T-t, x)$ we have $\varphi(0, x)=\phi(T, x)$. Then, similarly to (2.13) we may use (6.1) to determine a descent direction for $J$ by the following discretization for $i=0, \ldots, N_{x}$, and with $\varphi_{i}^{0}$ as discretization of $\varphi\left(0, x_{i}\right)$.

$$
\begin{align*}
& \varphi_{2 i}^{(1), 0}=\phi_{2 i}^{(1), N_{t}}=\frac{\eta_{2 i}^{(1), N_{t}}-\eta_{2 i}^{(2), N_{t}}}{2 a}-\left(y_{d}\right)_{2 i}, \varphi_{2 i+1}^{(1), 0}=\varphi_{2 i}^{(1), 0},  \tag{6.2a}\\
& \varphi_{2 i-1}^{(2), 0}=\phi_{2 i-1}^{(2), N_{t}}=-\left(\frac{\eta_{2 i-1}^{(1), N_{t}}-\eta_{2 i-1}^{(2), N_{t}}}{2 a}-\left(y_{d}\right)_{2 i-1}\right), \varphi_{2 i}^{(2), 0}=\varphi_{2 i-1}^{(2), 0} . \tag{6.2b}
\end{align*}
$$

In the present case we have $N\left(u_{0}\right)=N\left(\eta_{0}\right)=N_{x}$ since at each cell interface there is a discontinuity in the initial data $\eta_{0}$. Recall, that we have discontinuities in the first family only at odd indices and in the second family at even indices. We prescribe $\xi_{j}(T)$ for the $j$ th discontinuity at terminal time and propagate backwards in time by equation (2.10), i.e., $\xi_{j}(T)=\xi_{j}(0)$ for all discontinuities $j=0, \ldots, N_{x}$. Denote by $\Delta_{j} w(\cdot)=w\left(x_{j+1}\right)-w\left(x_{j}\right)$ and by $\bar{\Delta}_{j} w(\cdot)=\frac{1}{2}\left(w\left(x_{j+1}\right)+w\left(x_{j}\right)\right)$. Due to the choice of $\Delta t=\frac{\Delta x}{a}$ and $T=N_{t} \Delta t$, we have that a shock $j$ located initially at the cell interface $x_{j}(0)=x_{j+\frac{1}{2}}$ and moving at speed $a$ is located at time $T$ at $x_{j}(T)$ being again a cell interface located at $x_{i+\frac{1}{2}}=x_{j+\frac{1}{2}}+a T$ for some $i \in\left\{0, \ldots, N_{x}\right\}$. We denote by $i(j)$ this location. Note that for $j$ odd (even) the speed of the discontinuity is $-a(a)$ as in equation (2.25). The $j$ th discontinuity is by convention located initially at $x_{j}(0)=x_{j+\frac{1}{2}}$.

$$
\begin{align*}
\xi_{j} & =\xi_{j}(T)=\frac{1}{2}\left(\left(\bar{\Delta}_{i(j)} \eta^{(1), N_{t}}-\bar{\Delta}_{i(j)} \eta^{(2), N_{t}}-\bar{\Delta}_{i(j)} y_{d}\right)\right) \times  \tag{6.2c}\\
& \left(\widehat{\Delta}_{i(j)} \eta^{(1), N_{t}}-\widehat{\Delta}_{i(j)} \eta^{(2), N_{t}}\right) .
\end{align*}
$$

where $\bar{\Delta} w_{k}=\frac{1}{2 a}\left(w_{k+1}+w_{k}\right)$ and $\widehat{\Delta} w_{k}=w_{k+1}-w_{k}$. Note that the formula simplifies if this term is evaluated in original variables $y$ instead of $\eta$. However, in order to show the relation to equation (6.1) we write here the technical formulation.

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