# Godunov scheme and sampling technique for computing phase transitions in traffic flow modeling 

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#### Abstract

A new version of Godunov scheme is proposed in order to compute solutions of a traffic flow model with phase transitions. The scheme is based on a modified averaging strategy and a sampling procedure. Several numerical tests are shown to prove the validity of the method. The convergence of the algorithm is demonstrated numerically. We also give a higher order extension of the method in space and time.


Key words: Hyperbolic conservation laws, Continuous traffic models, Godunov scheme, Phase transitions, Sampling.

## 1 Introduction

We are interested in the numerical approximation of the solutions of a continuous traffic flow model taking into account the phenomenon of phase transitions between a free and a congested environment. The model we consider has been introduced by Colombo [10]. It consists of a scalar conservation law describing the free flow, and of a $2 \times 2$ system of conservation laws when the flow is congested. The coupling is achieved by introducing phase transitions between the free and the congested phase. We recall that, from the analytical point of view, the model is well posed for all initial data with bounded total variation [11].

Other traffic flow models with phase transitions have been considered in the literature since the 60ties, in order to explain empirical flow-density relations (see Helbing [16, Section II] for description of the features recovered by a detailed analysis of the fundamental diagram). In particular, we refer the reader to the scalar model of Drake, Schofer and May [13]. Another model has been introduced recently by the second author [14].

[^0]Historically, one of the first continuous models introduced to describe traffic flow is the well known Lighthill-Whitham [20] and Richards [21] (LWR) model, which reads

$$
\begin{equation*}
\partial_{t} \rho+\partial_{x}[\rho v(\rho)]=0, \tag{1}
\end{equation*}
$$

where $\rho \in[0, R]$ is the mean traffic density, and $v(\rho)$, the mean traffic velocity, is a given nonincreasing function, non-negative for $\rho$ between 0 and the positive maximal density $R$, which corresponds to a traffic jam. This scalar model expresses conservation of the number of cars, and relies on the assumption that the car speed depends only on the density (more complex closure relations between speed and density, involving the density gradient, can be assumed, see [2] and references therein). This phenomenological relation is valid in steady state conditions, and is not realistic in more complicated situations. In particular, as shown in Figure 1, the corresponding fundamental diagram in the $(\rho, \rho v)$-plane does not qualitatively match experimental data at high densities.



Figure 1: Left: standard flow for the LWR model. Right: experimental data, taken from [18]; here $q$ denotes the flux $\rho v$.

The diagram above suggests that a good traffic flow model should exhibit two qualitative different behaviors:

- for low densities, the flow is free and essentially analogous to the LWR model;
- at high densities the flow is congested and covers a 2 -dimensional domain in the fundamental diagram; a "second order" model seems more appropriate to describe this dynamic.

From a numerical point of view, the presence of phase transitions make standard numerical schemes unuseful. For example, it is easy to see that the classical Godunov method is not applicable due to the lack of convexity of the whole model phase space. Indeed, the latter turns out to be a disconnected set in $\mathbb{R}^{2}$, made of two connected components associated with the free and the congested domains, respectively. In presence of phase transitions, the projection step taking place in the classical Godunov method can then give values which are not in the domain. This necessarily stops the procedure. We are thus led to propose a new version of Godunov method, based on a modified averaging strategy and a sampling procedure. More precisely, we modify the mesh cells following the phase boundaries, so that the projection involves only values belonging to the same phase. In order to come back to the original cells, we complete the projection step with a Glimmtype sampling technique.

This scheme is essentially first order accurate, and hence introduces a considerable dissipation away from phase transitions. In order to improve accuracy, we also present an extension of the method to second-order accuracy in space and time, which is $L^{1}$-stable in space.

We remark that the techniques presented here apply also to models in $[13,14]$.
The averaging procedure on modified cells that we introduce has first been used (up to our knowledge) in [25] but in a different context and a slightly different form. However, the idea of going back to the initial cells by means of a sampling procedure is new and allows us to avoid dealing with moving meshes (as in [25]). It has been motivated by recent works proposed by the first author for approximating nonclassical solutions arising in certain nonlinear hyperbolic equations (see [4], [5] and the references therein), and very recently by Chalons and Coquel in [6] for computing sharp discrete shock profiles. Let us underline that the model studied here differs a lot from the one addressed by [25], since we are coupling systems of different dimensions. Moreover, we describe a higher order strategy adapted to our model.

The paper is organized as follows. In Section 2 we present the model, whose Riemann solver is described in Section 3. Section 4 is devoted to the description of the modified Godunov scheme and its higher order extention, which are tested in Section 5. Finally, Section 6 is devoted to the computation of the conservation errors.

## 2 Governing equations

We deal with a mathematical model introduced by Colombo [10] that aims at describing some experimental features observed in practice, like a discontinuous flow-density relation, "synchronized flow" and "wide moving jams" [18]. In particular, the associated system takes into account the phenomenon of phase transitions between a free flow and a congested flow.

For the free flow, the well known LWR scalar conservation law is used

$$
\begin{equation*}
\partial_{t} \rho+\partial_{x} \rho v_{f}(\rho)=0, \quad x \in \mathbb{R}, \quad t>0 \tag{2}
\end{equation*}
$$

where $\rho$ is the car density and $\rho \rightarrow v_{f}(\rho)$ is the speed function given by

$$
\begin{equation*}
v_{f}(\rho)=V\left(1-\frac{\rho}{R}\right) \tag{3}
\end{equation*}
$$

When the traffic is congested, the model reads as a $2 \times 2$ system of conservation laws [9]:

$$
\left\{\begin{array}{l}
\partial_{t} \rho+\partial_{x}\left(\rho v_{c}(\rho, q)\right)=0,  \tag{4}\\
\partial_{t} q+\partial_{x}\left((q-Q) v_{c}(\rho, q)\right)=0,
\end{array} \quad x \in \mathbb{R}, t>0\right.
$$

where $\rho$ still denotes the car density, $q$ is the weighted linear momentum and $Q$ is a parameter depending on the considered road and is related to the phenomenon of wide moving jams. The speed law now writes

$$
\begin{equation*}
v_{c}(\rho, q)=\left(1-\frac{\rho}{R}\right) \frac{q}{\rho} \tag{5}
\end{equation*}
$$

In (3) and (5), $R$ and $V$ are constants for the maximal possible car density and the maximal possible speed respectively.

It is worth noticing that one may provide the free system with a natural value for $q$ when setting $q=$ $\rho V$. Indeed $v_{f}(\rho)=v_{c}(\rho, \rho V)$, so that the associated conservation law is obtained by multiplying (2) by $V$ and writes

$$
\partial_{t} q+\partial_{x}\left(q v_{f}(\rho)\right)=0, \quad q=\rho V
$$

or equivalently

$$
\partial_{t} q+\partial_{x}\left(q v_{c}(\rho, q)\right)=0, \quad q=\rho V
$$

Recall that for the congested phase, $q$ evolves according to

$$
\partial_{t} q+\partial_{x}\left((q-Q) v_{c}(\rho, q)\right)=0
$$

This means that $q$ is not conserved across a phase transition. On the contrary, $\rho$ is conserved, as we will see in Section 3.

As pointed out in [12], empirical evidence suggests that queues cannot form spontaneously in free flow traffic for no apparent reason. Hence it is reasonable to assume that if the initial data are entirely in the free (resp. congested) phase, then the solution will belong to the free (resp. congested) phase for all time. Thus the phase space $\Omega_{f}$ (resp. $\Omega_{c}$ ) for the free system (resp. the congested system) is chosen to be an invariant sets for (2) (resp. (4)). (We refer the reader to [17] for definition and properties of invariant sets.) The domains are chosen to be

$$
\begin{align*}
\Omega_{f} & =\left\{(\rho, q) \in[0, R] \times\left[0,+\infty\left[, v_{f}(\rho) \geq V_{f}, q=\rho V\right\}\right.\right. \\
& =\left\{(\rho, q) \in[0, R] \times\left[0,+\infty\left[\frac{\rho}{R} \leq 1-\frac{V_{f}}{V}, q=\rho V\right\}\right.\right. \tag{6}
\end{align*}
$$

and

$$
\begin{equation*}
\Omega_{c}=\left\{(\rho, q) \in[0, R] \times\left[0,+\infty\left[, 0 \leq v_{c}(\rho, q) \leq V_{c}, \frac{Q^{-}-Q}{R} \leq \frac{q-Q}{\rho} \leq \frac{Q^{+}-Q}{R}\right\}\right.\right. \tag{7}
\end{equation*}
$$

where $V_{f}>V_{c}$ are thresholds speeds such that above $V_{f}$ the flow is free and below $V_{c}$ the flow is congested. Notice that the strict inequality $V_{f}>V_{c}$ is necessary to have uniqueness of solutions. In addition, $\left.Q^{-} \in\right] 0, Q\left[\right.$ and $\left.Q^{+} \in\right] Q,+\infty[$ depend from the environmental conditions and determine the width of the congested phase.

We introduce the following shorten form

$$
\begin{equation*}
\partial_{t} \mathbf{u}+\partial_{x} \mathbf{f}(\mathbf{u})=0, \quad \mathbf{u} \in \Omega=\Omega_{f} \cup \Omega_{c} \tag{8}
\end{equation*}
$$

for the model of phase transitions under consideration, with

$$
\left\{\begin{array}{lll}
\mathbf{u}=(\rho, q) & \text { and } \quad \mathbf{f}(\mathbf{u})=\left(\rho v_{f}(\rho), q v_{f}(\rho)\right), & \text { if } \quad(\rho, q) \in \Omega_{f} \\
\mathbf{u}=(\rho, q) & \text { and } \quad \mathbf{f}(\mathbf{u})=\left(\rho v_{c}(\rho, q),(q-Q) v_{c}(\rho, q)\right), & \text { if } \quad(\rho, q) \in \Omega_{c}
\end{array}\right.
$$

In the forthcoming developments, it is important to keep in mind that $\mathbf{u}$ and $\mathbf{f}(\mathbf{u})$ have not the same meaning in the free phase and in the congested phase. Moreover, the domain $\Omega$ is not connected, hence it is not convex (see Figure 2), left, in Section 3).
To conclude this section, let us underline that (8) is now supplemented, as is customary, with a given value of the solution at time $t=0$. More precisely, we assume as given an initial data $\mathbf{u}_{0} \in \Omega$ and we set

$$
\begin{equation*}
\mathbf{u}(., t=0)=\mathbf{u}_{0} \tag{9}
\end{equation*}
$$

## 3 About the Riemann problem

We recall in this section the description of the Riemann solver for (8)-(9), i.e. the self-similar solution of the Cauchy problem

$$
\left\{\begin{array}{l}
\partial_{t} \mathbf{u}+\partial_{x} \mathbf{f}(\mathbf{u})=0,  \tag{10}\\
\mathbf{u}_{0}(x)= \begin{cases}\mathbf{u}^{l}, & \text { if } x<0, \\
\mathbf{u}^{r}, & \text { if } x>0 .\end{cases}
\end{array}\right.
$$

If initial data $\mathbf{u}^{l}, \mathbf{u}^{r}$ are in the same phase, standard Lax solutions to the corresponding Riemann problem can be considered. Otherwise, following [10], admissible solutions are defined as follows.

Definition 1 If $\mathbf{u}^{l} \in \Omega_{f}$ and $\mathbf{u}^{r} \in \Omega_{c}$, then an admissible solution to (10) is a self-similar function $\mathbf{u}: \mathbb{R} \times\left[0,+\infty\left[\mapsto \Omega_{f} \cup \Omega_{c}\right.\right.$ such that, for some $\Lambda \in \mathbb{R}$, we have:

1. $\mathbf{u}(]-\infty, \Lambda t[, t) \subseteq \Omega_{f}$ and $\mathbf{u}(] \Lambda t,+\infty[, t) \subseteq \Omega_{c}$;
2. with clear notations, the functions

$$
\begin{align*}
& \mathbf{u}^{-}(x, t)= \begin{cases}\mathbf{u}(x, t) & \text { if } x<\Lambda t, \\
\mathbf{u}(\Lambda t-, t) & \text { if } x>\Lambda t,\end{cases}  \tag{11}\\
& \mathbf{u}^{+}(x, t)= \begin{cases}\mathbf{u}(\Lambda t+, t) & \text { if } x<\Lambda t, \\
\mathbf{u}(x, t) & \text { if } x>\Lambda t,\end{cases} \tag{12}
\end{align*}
$$

are Lax entropy solutions to corresponding Riemann problems for (2) and (4), respectively;
3. the Rankine-Hugoniot condition

$$
\rho(\Lambda t+, t) v_{c}(\mathbf{u}(\Lambda t+, t))-\rho(\Lambda t-, t) v_{f}(\rho(\Lambda t+, t))=\Lambda(\rho(\Lambda t+, t)-\rho(\Lambda t-, t))
$$

holds for all $t>0$.

If $\mathbf{u}^{l} \in \Omega_{c}$ and $\mathbf{u}^{r} \in \Omega_{f}$, the conditions are obtained by exchanging the roles of $\Omega_{f}, \Omega_{c}$ and $v_{f}, v_{c}$.
Notice that condition 3 above ensures that the total number of car is conserved across phase transitions.

Definition 1 does not assure uniqueness. We are then led to introduce the notion of consistency [10].

Definition 2 Let $\mathcal{R}:\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right) \mapsto \mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right)$ denote a Riemann solver, i.e. $x \mapsto \mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right)(x)$ is the solution of (10) computed at time $t=1 . \mathcal{R}$ is consistent if the following two conditions hold for all $\mathbf{u}^{l}, \mathbf{u}^{m}, \mathbf{u}^{r} \in \Omega_{f} \cup \Omega_{c}$, and $\bar{x} \in \mathbb{R}$ :

$$
\begin{align*}
& \left.\begin{array}{ll}
\mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right)(\bar{x})= & =\mathbf{u}^{m} \\
\mathcal{R}\left(\mathbf{u}^{m}, \mathbf{u}^{r}\right)(\bar{x}) & =\mathbf{u}^{m}
\end{array}\right\} \Rightarrow \mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right)= \begin{cases}\mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right), & \text { if } x<\bar{x}, \\
\mathcal{R}\left(\mathbf{u}^{m}, \mathbf{u}^{r}\right), & \text { if } x \geq \bar{x},\end{cases}  \tag{C1}\\
& \mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right)(\bar{x})=\mathbf{u}^{m} \Rightarrow \begin{cases}\mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right)= \begin{cases}\mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right), & \text { if } x \leq \bar{x}, \\
\mathbf{u}^{m}, & \text { if } x>\bar{x},\end{cases} \\
\mathcal{R}\left(\mathbf{u}^{m}, \mathbf{u}^{r}\right)= \begin{cases}\mathbf{u}^{m}, & \text { if } x<\bar{x}, \\
\mathcal{R}\left(\mathbf{u}^{l}, \mathbf{u}^{r}\right), & \text { if } x \geq \bar{x} .\end{cases} \end{cases} \tag{C2}
\end{align*}
$$





Figure 2: Notation used in the paper.

Essentially, (C1) states that whenever two solutions to two Riemann problems can be placed side by side, then their juxtaposition is again a solution to a Riemann problem. (C2) is the vice-versa.

We are now ready to construct the Riemann solver. We recall first the basic features of models (2) and (4). In the free phase the characteristic speed is $\lambda(\rho)=V(1-2 \rho / R)$, while the informations on system (4) are collected in the following table:

$$
\begin{array}{ll}
r_{1}(\rho, q)=\left[\begin{array}{c}
\rho \\
q-Q
\end{array}\right], & r_{2}(\rho, q)=\left[\begin{array}{c}
R-\rho \\
\frac{R}{\rho} q
\end{array}\right] \\
\lambda_{1}(\rho, q)=\left(\frac{2}{R}-\frac{1}{\rho}\right) \cdot(Q-q)-\frac{Q}{R}, & \lambda_{2}(\rho, q)=v_{c}(\rho, q) \\
\nabla \lambda_{1} \cdot r_{1}(\rho, q)=2 \frac{Q-q}{R}, & \nabla \lambda_{2} \cdot r_{2}(\rho, q)=0 \\
\mathcal{L}_{1}\left(\rho ; \rho_{o}, q_{o}\right)=Q+\frac{q_{o}-Q}{\rho_{o}} \rho, & \mathcal{L}_{2}\left(\rho ; \rho_{o}, q_{o}\right)=\frac{\rho}{\rho_{o}} \frac{R-\rho_{o}}{R-\rho} q_{o},  \tag{14}\\
w_{1}(\rho, q)=v_{c}(\rho, q), & w_{2}(\rho, q)=\frac{q-Q}{\rho}
\end{array}
$$

where $r_{i}$ is the $i$-th right eigenvector, $\lambda_{i}$ the corresponding eigenvalue, $\mathcal{L}_{i}$ is the $i$-Lax curve and $w_{i}$ is the $i$-Riemann invariant. Shock and rarefaction curves coincide, hence system (4) belongs to Temple class [22].

Using Riemann coordinates $\left(w_{1}, w_{2}\right), \Omega_{c}=\left[0, V_{c}\right] \times\left[W_{2}^{-}, W_{2}^{+}\right]$. For $(\rho, q) \in \Omega_{f}$, we extend the corresponding Riemann coordinates $\left(w_{1}, w_{2}\right)$ as follows. Let $\tilde{\mathbf{u}}=(\tilde{\rho}, \tilde{\rho} V)$ be the point in $\Omega_{f}$ defined by $\tilde{\rho}=Q /\left(V-W_{2}^{-}\right)$. Define

$$
w_{1}=V_{f} \quad \text { and } \quad w_{2}= \begin{cases}V-Q / \rho & \text { if } \rho \geq \tilde{\rho}  \tag{15}\\ v_{f}(\tilde{\rho})-v_{f}(\rho)+V-Q / \tilde{\rho} & \text { if } \rho<\tilde{\rho}\end{cases}
$$

so that, in the Riemann coordinates, $\Omega_{f}=\left\{V_{f}\right\} \times\left[W_{o}, W_{2}^{+}\right]$, see Figure 2.

All possible cases are listed below. Figures 3-5 are plotted using the parameters values given in Section 5.
(A) The data in (10) are in the same phase, i.e. they are either both in $\Omega_{f}$ or both in $\Omega_{c}$. Then the solution is the standard Lax entropy solution to (2), resp. (4), and no phase boundary is present.
(B) $\mathbf{u}^{l} \in \Omega_{c}$ and $\mathbf{u}^{r} \in \Omega_{f}$ (as in Figure 3 ). We consider the points $\mathbf{u}^{c} \in \Omega_{c}$ and $\mathbf{u}^{m} \in \Omega_{f}$ implicitly defined by

$$
\begin{aligned}
\left(1-\frac{\rho^{c}}{R}\right)\left(Q+w_{2}\left(\mathbf{u}^{l}\right) \rho^{c}\right) & =\rho^{c} V_{c} \\
\left(1-\frac{\rho^{m}}{R}\right)\left(Q+w_{2}\left(\mathbf{u}^{l}\right) \rho^{m}\right) & =\rho^{m} V\left(1-\frac{\rho^{m}}{R}\right)
\end{aligned}
$$

If $w_{2}\left(\mathbf{u}^{l}\right)>0$, the solution is made of a 1-rarefaction from $\mathbf{u}^{l}$ to $\mathbf{u}^{c}$, a phase transition from $\mathbf{u}^{c}$ to $\mathbf{u}^{m}$ and a Lax wave from $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: Test E). If $w_{2}\left(\mathbf{u}^{l}\right) \leq 0$, we have a shock-like phase transition from $\mathbf{u}^{l}$ to $\mathbf{u}^{m}$ and a Lax wave from $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: TestF).


Figure 3: Hugoniot locuses in the phase plane $(\rho, \rho v)$ for case (B).
(C) $\mathbf{u}^{l} \in \Omega_{f}$ and $\mathbf{u}^{r} \in \Omega_{c}$ with $w_{2}\left(\mathbf{u}^{l}\right) \in\left[W_{2}^{-}, W_{2}^{+}\right]$(Figure 4). Consider the points $\mathbf{u}^{c}$ and $\mathbf{u}^{m} \in \Omega_{c}$ implicitly defined by

$$
\begin{aligned}
\left(1-\frac{\rho^{c}}{R}\right)\left(Q+w_{2}\left(\mathbf{u}^{l}\right) \rho^{c}\right) & =\rho^{c} V_{c} \\
\left(1-\frac{\rho^{m}}{R}\right)\left(Q+w_{2}\left(\mathbf{u}^{l}\right) \rho^{m}\right) & =\rho^{m} w_{1}\left(\mathbf{u}^{r}\right)
\end{aligned}
$$

If $w_{2}\left(\mathbf{u}^{l}\right)>0$, the solution is made of a shock-like phase transition from $\mathbf{u}^{l}$ to $\mathbf{u}^{m}$ and a 2 -contact discontinuity $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: Test G). If $w_{2}\left(\mathbf{u}^{l}\right) \leq 0$, the solution displays a phase transition from $\mathbf{u}^{l}$ to $\mathbf{u}^{c}$, a 2-rarefaction from $\mathbf{u}^{c}$ to $\mathbf{u}^{m}$ and a 2-contact discontinuity $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: Test H).
(D) $\mathbf{u}^{l} \in \Omega_{f}$ with $w_{2}\left(\mathbf{u}^{l}\right)<W_{2}^{-}$and $\mathbf{u}^{r} \in \Omega_{c}$ (see Figure 5). Let $\mathbf{u}^{m} \in \Omega_{c}$ be the point on the lower boundary of $\Omega_{c}$ implicitly defined by

$$
\left(1-\frac{\rho^{m}}{R}\right)\left(Q+W_{2}^{-} \rho^{m}\right)=\rho^{m} w_{1}\left(\mathbf{u}^{r}\right)
$$

and consider the speed of the phase boundary joining $\mathbf{u}^{l} \in \Omega_{f}$ to $\mathbf{u}^{m} \in \Omega_{c}$

$$
\Lambda\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right)=\frac{\rho^{l} v_{f}\left(\rho^{l}\right)-\rho^{m} w_{1}\left(\mathbf{u}^{r}\right)}{\rho^{l}-\rho^{m}} .
$$



Figure 4: Hugoniot locuses in the phase plane $(\rho, \rho v)$ for case (C).

Let $\mathbf{U}_{c}=\left(R_{c}, Q_{c}\right) \in \Omega_{c}$ be the point whose Riemann coordinates are $\left(V_{c}, W_{2}^{-}\right)$. If $\lambda_{1}\left(\mathbf{U}_{c}\right) \geq$ $\Lambda\left(\mathbf{u}^{l}, \mathbf{U}_{c}\right)$, the solution is a phase transition from $\mathbf{u}^{l}$ to $\mathbf{U}_{c}$, a 1-rarefaction from $\mathbf{U}_{c}$ to $\mathbf{u}^{m}$ and a 2 -contact discontinuity from $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$. Otherwise:

- If $\lambda_{1}\left(\mathbf{u}^{m}\right) \leq \Lambda\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right)$, the solution is a phase transition from $\mathbf{u}^{l}$ to $\mathbf{u}^{m}$ followed by a 2-contact discontinuity from $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: Test J).
- If $\lambda_{1}\left(\mathbf{u}^{m}\right)>\Lambda\left(\mathbf{u}^{l}, \mathbf{u}^{m}\right)$, let $\mathbf{u}^{c}=\left(\rho^{c}, q^{c}\right) \in \Omega_{c}$ implicitly defined by

$$
\lambda_{1}\left(\mathbf{u}^{c}\right)=\Lambda\left(\mathbf{u}^{l}, \mathbf{u}^{c}\right),
$$

i.e. $\rho^{c}$ is the bigger root of the equation

$$
\left(Q-Q^{-}\right) \rho^{2}-2 \rho^{l}\left(Q-Q^{-}\right) \rho+R^{2}\left(\rho^{l} v_{f}\left(\rho^{l}\right)-Q\right)+\rho^{l} R\left(2 Q-Q^{-}\right)=0
$$

and $q^{c}=Q-\rho^{c}\left(Q-Q^{-}\right) / R$. Then the solution shows a phase transition from $\mathbf{u}^{l}$ to $\mathbf{u}^{c}$, an attached 1-rarefaction from $\mathbf{u}^{c}$ to $\mathbf{u}^{m}$ and a 2-contact discontinuity from $\mathbf{u}^{m}$ to $\mathbf{u}^{r}$ (Section 5: Test I).

## 4 Numerical schemes

Let us first introduce a space step $\Delta x$ and a time step $\Delta t$, both assumed to be constant for simplicity in the forthcoming developments. We set $\nu=\Delta t / \Delta x$. Then, we define the mesh interfaces $x_{j+1 / 2}=j \Delta x$ for $j \in \mathbb{Z}$ and the intermediate times $t^{n}=n \Delta t$ for $n \in \mathbb{N}$, and we seek at each time $t^{n}$ an approximation $\mathbf{u}_{j}^{n}$ of the solution of (8)-(9) on the interval $\left[x_{j-1 / 2}, x_{j+1 / 2}\right), j \in \mathbb{Z}$. Therefore, a piecewise constant approximated solution $x \rightarrow \mathbf{u}_{\nu}\left(x, t^{n}\right)$ of the solution $\mathbf{u}$ is given by

$$
\mathbf{u}_{\nu}\left(x, t^{n}\right)=\mathbf{u}_{j}^{n} \text { for all } x \in C_{j}=\left[x_{j-1 / 2} ; x_{j+1 / 2}\right), \quad j \in \mathbb{Z}, \quad n \in \mathbb{N}
$$

When $n=0$, we set $x_{j}=0.5 \cdot\left(x_{j-1 / 2}+x_{j+1 / 2}\right)$ and

$$
\mathbf{u}_{j}^{0}=\mathbf{u}_{0}\left(x_{j}\right), \text { for all } j \in \mathbb{Z}
$$



Figure 5: Hugoniot locuses in the phase plane $(\rho, \rho v)$ for case (D).

Note that the usual $L^{2}$-projection is not adapted in the present context since, depending on initial data, it could artificially introduce unphysical states which are not in the phase space at time $t=0$ (recall that $\Omega=\Omega_{f} \cup \Omega_{c}$ is not convex).

Assuming as given a sequence $\left(\mathbf{u}_{j}^{n}\right)_{j \in \mathbb{Z}}$ at time $t^{n}$, it is now a question of proposing a way of advancing it to the next time level $t^{n+1}$. Two manners of tackling that are now presented. We begin with the celebrated Godunov scheme in its classical form and show that it is not relevant in the present setting of possible phase transitions. We are thus led to present a new version of this scheme based on a modified averaging strategy coupled with a sampling procedure. We also show an extension of this last method to second-order accuracy.

### 4.1 Failure of the classical Godunov scheme

As it is well-known, the Godunov scheme is composed of two steps : a first step in which initial data evolves in time according to the PDE model under consideration, and a second step of projection onto piecewise constant functions. Let us give more details.

## Step 1: Evolution in time

In this first step, one solves the following Cauchy problem

$$
\left\{\begin{array}{l}
\partial_{t} \mathbf{v}+\partial_{x} \mathbf{f}(\mathbf{v})=0, \quad x \in \mathbb{R},  \tag{16}\\
\mathbf{v}(x, 0)=\mathbf{u}_{\nu}\left(x, t^{n}\right)
\end{array}\right.
$$

for times $t \in[0, \Delta t]$. Recall that $x \rightarrow \mathbf{u}_{\nu}\left(x, t^{n}\right)$ is piecewise constant. Then, under the usual CFL restriction

$$
\begin{equation*}
\frac{\Delta t}{\Delta x} \max _{\mathbf{v}}\left\{\left|\lambda_{i}(\mathbf{v})\right|, i=1 \text { if } \mathbf{v} \in \Omega_{f}, i=1,2 \text { if } \mathbf{v} \in \Omega_{c}\right\} \leq \frac{1}{2} \tag{17}
\end{equation*}
$$

for all the $\mathbf{v}$ under consideration, the solution of (16) is known by gluing together the solutions of the Riemann problems set at each interface. More precisely

$$
\begin{equation*}
\mathbf{v}(x, t)=\mathbf{v}_{\mathbf{r}}\left(\frac{x-x_{j+1 / 2}}{t} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right) \text { for all }(x, t) \in\left[x_{j}, x_{j+1}\right] \times[0, \Delta t] \tag{18}
\end{equation*}
$$

where $(x, t) \rightarrow \mathbf{v}_{\mathbf{r}}\left(\frac{x}{t} ; \mathbf{v}^{l}, \mathbf{v}^{r}\right)$ denotes the self-similar solution of the Riemann problem

$$
\left\{\begin{array}{l}
\partial_{t} \mathbf{v}+\partial_{x} \mathbf{f}(\mathbf{v})=0, \\
\mathbf{v}(x, 0)=\left\{\begin{array}{lll}
\mathbf{v}^{l} & \text { if } & x<0, \\
\mathbf{v}^{r} & \text { if } & x>0,
\end{array}\right.
\end{array}\right.
$$

whatever $\mathbf{v}^{l}$ and $\mathbf{v}^{r}$ are in the phase space $\Omega_{f} \cup \Omega_{c}$ (see also Section 3 above).

Step 2 : Projection ( $t^{n} \rightarrow t^{n+1}$ )
Aim of this second step is to get a piecewise constant approximated solution on each cell $\mathcal{C}_{j}$ at time $t^{n+1}$. This may be simply done by averaging the solution $x \rightarrow \mathbf{v}(x, \Delta t)$ given by (18), as expressed by the following update formula:

$$
\begin{equation*}
\mathbf{u}_{j}^{n+1}=\frac{1}{\Delta x} \int_{x_{j-1 / 2}}^{x_{j+1 / 2}} \mathbf{v}(x, \Delta t) d t, \quad j \in \mathbb{Z} \tag{19}
\end{equation*}
$$

Actually, one can provide an even simpler formula for $\mathbf{u}_{j}^{n+1}$ by integrating equation (16) over the element $E=(a b c d)$ defined by $\left[x_{j-1 / 2}, x_{j+1 / 2}\right] \times[0, \Delta t]$ and represented on Figure 6. ¿From Green's


Figure 6: Averaging element in the classical Godunov method
theorem and using (16) and (19), we get

$$
\begin{aligned}
0 & =\iint_{E} \partial_{t} \mathbf{v}+\partial_{x} \mathbf{f}(\mathbf{v}) d x d t \\
& =\Delta x \mathbf{u}_{j}^{n+1}-\Delta x \mathbf{u}_{j}^{n}+\int_{0}^{\Delta t} \mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right) d t-\int_{0}^{\Delta t} \mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j-1}^{n}, \mathbf{u}_{j}^{n}\right)\right) d t
\end{aligned}
$$

where classical notations have been used for the traces at $0^{-}$and $0^{+}$of the Riemann solutions under consideration. Introducing the numerical fluxes

$$
\mathbf{f}_{j+1 / 2}^{n, \pm}=\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{ \pm} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right) \text { for all } j \in \mathbb{Z}
$$

we immediately obtain

$$
\begin{equation*}
\mathbf{u}_{j}^{n+1}=\mathbf{u}_{j}^{n}-\frac{\Delta t}{\Delta x}\left(\mathbf{f}_{j+1 / 2}^{n,-}-\mathbf{f}_{j-1 / 2}^{n,+}\right) \text { for all } j \in \mathbb{Z} \tag{20}
\end{equation*}
$$

To conclude the description of the classical Godunov scheme, let us observe that if $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$ belong to the same phase (free or congested), the same holds for the Riemann solution $(x, t) \rightarrow$ $\mathbf{v}_{\mathbf{r}}\left(\frac{x}{t} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ so that $\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ and $\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ are actually in the same phase. As a consequence, we necessarily have for all $j \in \mathbb{Z}$ :

$$
\begin{equation*}
\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)=\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right) . \tag{21}
\end{equation*}
$$

That is indeed obvious if $\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)=\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ while in the opposite case, the associated discontinuity is necessarily stationary and (21) follows from Rankine-Hugoniot conditions. In other words, the method is conservative at the corresponding interface $j+1 / 2$.
Otherwise, if $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$ are not in the same phase, the states $\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ and $\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ may well be distinct and in two different phases. In such a situation, only the first component of $\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)$ and $\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0^{+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)$ associated with the conservation of the mass in both the free and the congested system are equal. Indeed, recall that only the Rankine-Hugoniot relation coming from the mass conservation applies across a phase transition (see Definition 1 in Section 3 above).

## What is wrong with this method.

Of course, the computation of the traces at $0^{-}$and $0^{+}$of the Riemann solvers at each interface $x_{j+1 / 2}$ does not represent a problem, since these solutions are known from Section 3 (see also [10]). Actually, the failure of this strategy is due to the lack of convexity of the domain $\Omega_{f} \cup \Omega_{c}$ in the ( $\rho, q$ )-plane (see Figure 2) and to the possible presence of phase transitions in the Riemann solutions. In this case, the state $\mathbf{u}_{j}^{n+1}$ resulting from the averaging procedure (19) can be outside $\Omega_{f} \cup \Omega_{c}$ for some $j \in \mathbb{Z}$ (even if the solution $\mathbf{u}_{\nu}\left(., t^{n}\right)$ belongs to the domain). This means that at the next time step, the Riemann solutions are generally not known (even not defined actually) at each interface $x_{j+1 / 2}$, so that the classical Godunov method stops. As already discussed by Zhong, Hou and LeFloch in [25], and even earlier by Abgrall [1] for more simple and more homogeneous problems, the problem we underline here is clearly located in the projection step of the Godunov method (Step 2). In order to know, everywhere in the domain and at each time step, in which phase the flow is, we now propose a modification of this second step.

### 4.2 A new version of the Godunov scheme

In this section, we present a new version of the Godunov scheme that will turn out to be more adapted to the model under consideration. We keep unchanged the first step of the method as it is described in the previous section. Indeed, let us recall that the failure of the classical Godunov method comes from the projection strategy proposed in the second step. The difference thus lies in the corresponding averaging procedure (19) that we are going to modify. The idea (see [25] in a different context) is to no longer average the solution $x \rightarrow \mathbf{v}(x, \Delta t)$ on the mesh cells $\mathcal{C}_{j}=\left[x_{j-1 / 2}, x_{j+1 / 2}\right.$ ), since they may contain states in different phases, but on (possibly) modified and non uniform cells that we will denote $\overline{\mathcal{C}}_{j}^{n}=\left[\bar{x}_{j-1 / 2}^{n}, \bar{x}_{j+1 / 2}^{n}\right)$. These are constructed to contain values belonging to a single phase. To this aim, the new cells are defined according to the position of the phase transitions. Then, a sampling strategy will allow us to recover a piecewise constant solution on the initial mesh cells $\mathcal{C}_{j}$ (see also [4], [5] in a different context).

Step 2 (Modified) : Projection ( $t^{n} \rightarrow t^{n+1}$ )
Let $\left(\sigma_{j+1 / 2}^{n}=\sigma\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)_{j \in \mathbb{Z}}$ be a sequence of characteristic speeds of propagation at interfaces
$\left(x_{j+1 / 2}\right)_{j \in \mathbb{Z}}$ such that :

- if $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$ are not in the same phase (free or congested), then $\sigma_{j+1 / 2}^{n}$ coincides with the speed of propagation of the phase transition in the Riemann solution $(x, t) \rightarrow \mathbf{v}_{\mathbf{r}}\left(\frac{x}{t} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$,
- if $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$ belong to the same phase, then $\sigma_{j+1 / 2}^{n}=0$. Then, assuming that for all $j \in \mathbb{Z}$ the interface $x_{j+1 / 2}$ moves at velocity $\sigma_{j+1 / 2}^{n}$ between times $t^{n}$ and $t^{n+1}=t^{n}+\Delta t$, we define the new interface $\bar{x}_{j+1 / 2}^{n}$ at time $t^{n+1}$ setting

$$
\begin{equation*}
\bar{x}_{j+1 / 2}^{n}=x_{j+1 / 2}+\sigma_{j+1 / 2}^{n} \Delta t, \quad j \in \mathbb{Z} \tag{22}
\end{equation*}
$$

We also introduce

$$
\overline{\Delta x}_{j}^{n}=\bar{x}_{j+1 / 2}^{n}-\bar{x}_{j-1 / 2}^{n}, \quad j \in \mathbb{Z}
$$

In particular and by definition of the sequence $\left(\sigma_{j+1 / 2}^{n}\right)_{j \in \mathbb{Z}}$, it is clear that on each modified cell $\overline{\mathcal{C}}_{j}^{n}=\left[\bar{x}_{j-1 / 2}^{n}, \bar{x}_{j+1 / 2}^{n}\right.$ ), the solution $x \rightarrow \mathbf{v}(x, \Delta t)$ given by (18) is fully either in the free phase or in the congested phase. Then, averaging this solution on cells $\overline{\mathcal{C}}_{j}^{n}$ provide us with a piecewise constant approximated solution $\overline{\mathbf{u}}_{\nu}\left(x, t^{n+1}\right)$ on a non uniform mesh defined by

$$
\overline{\mathbf{u}}_{\nu}\left(x, t^{n+1}\right)=\overline{\mathbf{u}}_{j}^{n+1} \text { for all } x \in \bar{C}_{j}^{n}, \quad j \in \mathbb{Z}, \quad n \in \mathbb{N}
$$

with

$$
\overline{\mathbf{u}}_{j}^{n+1}=\frac{1}{\overline{\Delta x}}{ }_{j}^{n} \int_{\bar{x}_{j-1 / 2}^{n}}^{\bar{x}_{j+1 / 2}^{n}} \mathbf{v}(x, \Delta t) d t, \quad j \in \mathbb{Z}
$$

Let us underline that by the definition of the modified cells, we actually know which phase every constant state of the solution $\overline{\mathbf{u}}_{\nu}\left(x, t^{n+1}\right)$ belongs to. In fact, both $\Omega_{f}$ and $\Omega_{c}$ are convex domains (and then are stable under the process of an $L^{2}$-projection). This is the relevant difference with the solution $\mathbf{u}_{\nu}\left(x, t^{n+1}\right)$ obtained in the classical Godunov method, and justify the new approach. Let us notice that the modified cells $\overline{\mathcal{C}}_{j}^{n}$ may be either smaller or larger than the original ones $\mathcal{C}_{j}$, depending on the signs of the velocities $\sigma_{j+1 / 2}^{n}, j \in \mathbb{Z}$ This is illustrated on Figures 7 and 8 .


Figure 7: A first example of averaging element in the modified Godunov method


Figure 8: A second example of averaging element in the modified Godunov method

Even in this case, a simpler formula is obtained for $\overline{\mathbf{u}}_{j}^{n+1}$ by integrating equation (16) over the element $\bar{E}=(a b c d)$ defined by :

$$
\bar{E}=\left\{(x, t) \text { such that } t \in[0, \Delta t] \text { and } x_{j-1 / 2}+\sigma_{j-1 / 2}^{n} t \leq x \leq x_{j+1 / 2}+\sigma_{j+1 / 2}^{n} t\right\}
$$

(see again Figures 7 and 8). Nevertheless, one has to be careful when applying Green's theorem owing to the fact that $E$ is no longer a rectangle. We get, using the same (usual) notations as in the previous section for denoting the traces of the Riemann solutions at given points :

$$
\begin{aligned}
0 & =\iint_{\bar{E}} \partial_{t} \mathbf{v}+\partial_{x} \mathbf{f}(\mathbf{v}) d x d t \\
& =\overline{\Delta x}_{j}^{n} \overline{\mathbf{u}}_{j}^{n+1}-\Delta x \mathbf{u}_{j}^{n} \\
& +\int_{0}^{\Delta t} \mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)-\sigma_{j+1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right) d t \\
& -\int_{0}^{\Delta t} \mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j-1 / 2}^{n,+} ; \mathbf{u}_{j-1}^{n}, \mathbf{u}_{j}^{n}\right)\right)-\sigma_{j-1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j-1 / 2}^{n,+} ; \mathbf{u}_{j-1}^{n}, \mathbf{u}_{j}^{n}\right) d t
\end{aligned}
$$

Introducing now the numerical fluxes

$$
\begin{equation*}
\overline{\mathbf{f}}_{j+1 / 2}^{n, \pm}=\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n, \pm} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)-\sigma_{j+1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n, \pm} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right) \text { for all } j \in \mathbb{Z} \tag{23}
\end{equation*}
$$

a condensed form, similar to (20), is obtained for $\overline{\mathbf{u}}_{j}^{n+1}$ :

$$
\begin{equation*}
\overline{\mathbf{u}}_{j}^{n+1}=\frac{\Delta x}{\overline{\Delta x}} \mathbf{u}_{j}^{n}-\frac{\Delta t}{\overline{\Delta x}}{ }_{j}^{n}\left(\overline{\mathbf{f}}_{j+1 / 2}^{n,-}-\overline{\mathbf{f}}_{j-1 / 2}^{n,+}\right) \text { for all } j \in \mathbb{Z} \tag{24}
\end{equation*}
$$

At this stage, let us notice that if $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$ are in the same phase, then $\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ and $\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ also does. The conservation property

$$
\begin{gather*}
\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)-\sigma_{j+1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)  \tag{25}\\
= \\
\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)-\sigma_{j+1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)
\end{gather*}
$$

then remains valid thanks to Rankine-Hugoniot conditions. Actually, note that in such a situation $\sigma_{j+1 / 2}^{n}=0$ by definition. Otherwise, if $\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,-} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ and $\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n,+} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ are not in the same phase, equality (25) makes sense only for the first component associated with the mass conservation.
Finally, we introduce the notation $\overline{\mathbf{f}}_{j+1 / 2}^{n, \pm}=\overline{\mathbf{f}}^{ \pm}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ for the interface fluxes, with of course

$$
\begin{equation*}
\overline{\mathbf{f}}^{ \pm}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)=\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n, \pm} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)-\sigma_{j+1 / 2}^{n} \mathbf{v}_{\mathbf{r}}\left(\sigma_{j+1 / 2}^{n, \pm} ; \mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right) \tag{26}
\end{equation*}
$$

Recall that $\sigma_{j+1 / 2}^{n}=\sigma\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ by definition.
In order to avoid having to deal with moving meshes, we complete the projection step defining a new approximation $\mathbf{u}_{j}^{n+1}$ of the solution at time $t^{n+1}$ on the (uniform) cells $\mathcal{C}_{j}, j \in \mathbb{Z}$. To this aim, for all $j \in \mathbb{Z}$, we propose to pick up randomly on the cell $\mathcal{C}_{j}$ a value between $\overline{\mathbf{u}}_{j-1}^{n+1}, \overline{\mathbf{u}}_{j}^{n+1}$ and $\overline{\mathbf{u}}_{j+1}^{n+1}$,
in agreement with their rate of presence in the cell. More precisely, given a well distributed random sequence $\left(a_{n}\right)$ within interval $(0,1)$, it amounts to set :

$$
\mathbf{u}_{j}^{n+1}=\left\{\begin{array}{lll}
\overline{\mathbf{u}}_{j-1}^{n+1} & \text { if } & a_{n+1} \in\left(0, \frac{\Delta t}{\Delta x} \max \left(\sigma_{j-1 / 2}^{n}, 0\right)\right),  \tag{27}\\
\overline{\mathbf{u}}_{j}^{n+1} & \text { if } & a_{n+1} \in\left[\frac{\Delta t}{\Delta x} \max \left(\sigma_{j-1 / 2}^{n}, 0\right), 1+\frac{\Delta t}{\Delta x} \min \left(\sigma_{j+1 / 2}^{n}, 0\right)\right), \\
\overline{\mathbf{u}}_{j+1}^{n+1} & \text { if } & a_{n+1} \in\left[1+\frac{\Delta t}{\Delta x} \min \left(\sigma_{j+1 / 2}^{n}, 0\right), 1\right),
\end{array}\right.
$$

for all $j \in \mathbb{Z}$.
Following a proposal by Collela [8], we consider van der Corput random sequence $\left(a_{n}\right)$ defined by

$$
a_{n}=\sum_{k=0}^{m} i_{k} 2^{-(k+1)},
$$

where $n=\sum_{k=0}^{m} i_{k} 2^{k}, i_{k}=0,1$, denotes the binary expansion of the integers $n=1,2, \ldots$. For the sake of illustration, we easily find that the first few elements of this sequence are

$$
\begin{array}{llll}
a_{1}=0.5, & a_{2}=0.25, & a_{3}=0.75, & a_{4}=0.125 \\
a_{5}=0.625, & a_{6}=0.375, & a_{7}=0.875, & a_{8}=0.0625
\end{array}
$$

This sequence is actually well distributed within the interval $(0,1)$. Moreover, one can prove for instance that $a_{i}<0.5$ for $i$ even and $a_{i}>0.5$ for $i$ odd. This well-known sequence is often favorite since, when used in the context of Glimm scheme, it leads to very good results in the smooth parts of the solutions (see for instance [8] and [7] for illustration). This concludes the description of the modified Godunov scheme.

We finish this section by emphasizing that due to the sampling procedure, the whole algorithm we propose (Step $1+$ Step 2 (Modified)) is not "strictly" conservative in the classical sense of finite volumes methods. However, we numerically demonstrate in the next sections that it is actually "weakly" conservative in the following sense : first, phase transitions propagates with the right speeds (given by Rankine-Hugoniot conditions) and then, conservation errors seems to tend to zero with the mesh size.

## Remark.

Of course, the random choice method (Glimm scheme) can be applied successfully to compute solutions of (8). Nevertheless, our method doesn't need to compute all the values in the Riemann solution, but only the values on both sides of the phase transition. Moreover, the algorithm coincides with the classical Godunov scheme, and hence it is conservative, away from phase transitions.

### 4.3 Higher order extension of the method in space and time

In this section, we describe a both space and time second-order extension of the new version of the Godunov scheme. Our strategy relies on the very popular MUSCL approach for the space accuracy and on a Runge-Kutta technique for the time accuracy. As usual, the second-order accuracy is obtained for smooth solutions only, even if better numerical results are also expected when discontinuities (or not smooth regions) are present. In our context, it is important to notice that smooth solutions exist but necessarily remain in a same phase (free or congested) since phase changes are always associated with discontinuities. As a consequence, the resulting procedure has
to be understood as second-order accurate away from phase transitions. This is the reason why we will focus on the first part of the projection step only, the sampling procedure being kept unchanged. This section is organized as follows. We first address the space accuracy and show how to obtain a MUSCL scheme which is stable in the $L^{1}$ sense. Then, we deal with the time accuracy and show how to apply a second-order Runge-Kutta technique.

## Accuracy in space

We begin by briefly recalling the MUSCL method for obtaining the second-order accuracy in space. For more details we refer the reader to [24], [15], [23] and the references therein, and also to a recent work by Berthon [3]. Assume that there exists a change of variables $\mathbf{u} \rightarrow \mathbf{U}=\varphi(\mathbf{u})$ from $\Omega$ onto some set $\Omega_{\mathbf{U}}$. The starting point of the method consists in replacing at each time $t^{n}$ and on each cell $C_{j}$ the constant values $\mathbf{u}_{j}^{n}$ by means of $\varphi$ and a linear reconstruction of $\mathbf{U}$. We set

$$
\left\{\begin{array}{l}
\mathbf{u}^{n}(x)=\varphi^{-1}\left(\mathbf{U}^{n}(x)\right) x \in C_{j}=\left[x_{j-1 / 2} ; x_{j+1 / 2}\right), \text { with }  \tag{28}\\
\mathbf{U}^{n}(x)=\mathbf{U}_{j}^{n}+s_{j}^{n} \frac{\left.x-x_{j}\right)}{\Delta x}, \quad \mathbf{U}_{j}^{n}=\varphi\left(\mathbf{u}_{j}^{n}\right), j \in \mathbb{Z} .
\end{array}\right.
$$

In (28), $x_{j}$ represents the center of the cell $C_{j}: x_{j}=\frac{1}{2}\left(x_{j-1 / 2}+x_{j+1 / 2}\right)$, and $s_{j}^{n}$ is the slope of the linear reconstruction. The choice of the reconstructed variable $\mathbf{U}$ generally depends on the system under consideration. In the present study, we will see below that in order to ensure the $L^{1}$-stability of the method, the reconstruction can be performed on the conservative variable for the free system ( $\mathbf{U}=\mathbf{u}$ ) while the Riemann variables $\mathbf{U}=\left(w_{1}, w_{2}\right)$ turn out to be more adapted for the congested phase.
We denote $\mathbf{u}_{j}^{n, \pm}$ and $\mathbf{U}_{j}^{n, \pm}$ the values at the edges $x=x_{j \pm 1 / 2}$ of $\mathbf{u}^{n}$ and $\mathbf{U}^{n}$ respectively :

$$
\left\{\begin{array}{l}
\mathbf{u}_{j}^{n, \pm}=\varphi^{-1}\left(\mathbf{U}_{j}^{n, \pm}\right),  \tag{29}\\
\mathbf{U}_{j}^{n, \pm}=\mathbf{U}_{j}^{n} \pm \Delta \mathbf{U}_{j}^{n} \quad \text { with } \quad \Delta \mathbf{U}_{j}^{n}=\frac{1}{2} s_{j}^{n} .
\end{array}\right.
$$

Then, following the basic principle of the MUSCL method, we propose to replace the couple $\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ with $\left(\mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)$ in the evaluation of the numerical fluxes $\overline{\mathbf{f}}_{j+1 / 2}^{n, \pm}$ at each interface $j+1 / 2$. More precisely, we now set $\overline{\mathbf{f}}_{j+1 / 2}^{n, \pm}=\overline{\mathbf{f}}^{ \pm}\left(\mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)$ instead of $\overline{\mathbf{f}}_{j+1 / 2}^{n, \pm}=\overline{\mathbf{f}}^{ \pm}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ in (26).

We draw now a particular attention to the choice of the reconstructed variable $\mathbf{U}$ and the slopes $s_{j}^{n}$. It is well-known that these have to be carefully determined for stability reasons. Once $\mathbf{U}$ is chosen, an usual choice for $s_{j}^{n}$, or equivalently for $\Delta \mathbf{U}_{j}^{n}$, is given by a slope-limiter procedure with for instance the so-called minmod limiter. It reads

$$
\begin{equation*}
\Delta \mathbf{U}_{j}^{n}=\frac{1}{2} \operatorname{minmod}\left(\mathbf{U}_{j+1}^{n}-\mathbf{U}_{j}^{n}, \mathbf{U}_{j}^{n}-\mathbf{U}_{j-1}^{n}\right) \tag{30}
\end{equation*}
$$

where the minmod function is defined by

$$
\operatorname{minmod}(a, b)=\left\{\begin{array}{ccc}
\operatorname{sign}(a) \min (|a|,|b|) & \text { if } & a b \geq 0, \\
0 & \text { otherwise }
\end{array}\right.
$$

for two scalar quantities $a$ and $b$ (sign denotes the sign function). In (30), minmod is applied component by component.

## Remark.

In the rest of this section, in order to avoid cumbersome notations, we will assume that the three states $\mathbf{U}_{j-1}^{n}, \mathbf{U}_{j}^{n}$ and $\mathbf{U}_{j+1}^{n}$ in (30) belong to the same phase. Otherwise, if $\mathbf{U}_{j-1}^{n}$ and/or $\mathbf{U}_{j+1}^{n}$ are not in the same phase of $\mathbf{U}_{j}^{n}$, they are replaced in practice with $\varphi\left(\mathbf{u}_{+}\left(\mathbf{u}_{j-1}^{n}, \mathbf{u}_{j}^{n}\right)\right)$ and/or $\varphi\left(\mathbf{u}_{-}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)$, where $\mathbf{u}_{ \pm}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)$ represent for all $j$ the values on both side of the phase transition in the Riemann solution associated with initial states $\mathbf{u}_{j}^{n}$ and $\mathbf{u}_{j+1}^{n}$. Then, by definition $\varphi\left(\mathbf{u}_{+}\left(\mathbf{u}_{j-1}^{n}, \mathbf{u}_{j}^{n}\right)\right), \mathbf{u}_{j}^{n}$ and $\varphi\left(\mathbf{u}_{-}\left(\mathbf{u}_{j}^{n}, \mathbf{u}_{j+1}^{n}\right)\right)$ belong to the same phase.

Our objective is to enforce the $L^{1}$-stability of the reconstruction procedure in the sense that we want $\mathbf{u}_{j}^{n, \pm}$ to necessarily belong to the phase space $\Omega$. This stability property will be obtained by a relevant choice of $\mathbf{U}$ and after a possible correction of the "guess" increment defined by (30).

Let us begin with the case $\mathbf{u}_{j}^{n} \in \Omega_{f}$ and consider a conservative reconstruction : $\mathbf{U}=\mathbf{u}$. Since $q$ always equals $\rho V$ in the free phase, the constraints $\mathbf{u}_{j}^{n, \pm} \in \Omega_{f}$ read

$$
\left\{\begin{array}{l}
0 \leq \rho_{j}^{n, \pm} \leq R \\
V_{f} \leq v_{f}\left(\rho_{j}^{n, \pm}\right)
\end{array}\right.
$$

which by definition of $v_{f}$ is first equivalent to

$$
\left\{\begin{array}{l}
0 \leq \rho_{j}^{n, \pm} \leq R \\
\rho_{j}^{n, \pm} \leq R\left(1-\frac{V_{f}}{V}\right)
\end{array}\right.
$$

and then to

$$
0 \leq \rho_{j}^{n, \pm} \leq R\left(1-\frac{V_{f}}{V}\right)
$$

by positivity of $V_{f}$ and $V$. But $\rho_{j}^{n, \pm}=\rho_{j}^{n} \pm \Delta \rho_{j}^{n}$ so that a straightforward transformation shows that these two conditions write as follows on the increment $\Delta \rho_{j}^{n}$ :

$$
\begin{equation*}
\left|\Delta \rho_{j}^{n}\right| \leq \min \left(\rho_{j}^{n}, R\left(1-\frac{V_{f}}{V}\right)-\rho_{j}^{n}\right) \tag{31}
\end{equation*}
$$

Then, choosing the "guess" increment (30) under the constraint (31) immediately leads to the following definition of $\Delta \rho_{j}^{n}$ :
$\Delta \rho_{j}^{n}=\operatorname{sign}(\Delta \rho) \min \left(\rho_{j}^{n}, R\left(1-\frac{V_{f}}{V}\right)-\rho_{j}^{n}, \Delta \rho\right), \quad$ with $\quad \Delta \rho=\frac{1}{2} \operatorname{minmod}\left(\rho_{j+1}^{n}-\rho_{j}^{n}, \rho_{j}^{n}-\rho_{j-1}^{n}\right)$.

We now turn to the case $\mathbf{u}_{j}^{n} \in \Omega_{c}$. The constraints $\mathbf{u}_{j}^{n, \pm} \in \Omega_{c}$ are equivalent to

$$
\left\{\begin{array}{l}
0 \leq w_{1}\left(\rho_{j}^{n, \pm}, q_{j}^{n, \pm}\right) \leq V_{c}  \tag{32}\\
W_{2}^{-} \leq w_{2}\left(\rho_{j}^{n, \pm}, q_{j}^{n, \pm}\right) \leq W_{2}^{+}
\end{array}\right.
$$

Due to the nonlinearity of the two functions $w_{1}$ and $w_{2}$, we propose to perform the reconstruction on these variables (the Riemann coordinates) and to set $\mathbf{U}=\left(w_{1}, w_{2}\right)(\rho, q)$. Then choosing the corresponding increments $\Delta\left(w_{1}\right)_{j}^{n}$ and $\Delta\left(w_{2}\right)_{j}^{n}$ according to (30) yields

$$
\left\{\begin{aligned}
\left(w_{1}\right)_{j}^{n, \pm} & =\left(w_{1}\right)_{j}^{n} \pm \frac{1}{2} \operatorname{minmod}\left(\left(w_{1}\right)_{j+1}^{n}-\left(w_{1}\right)_{j}^{n},\left(w_{1}\right)_{j}^{n}-\left(w_{1}\right)_{j-1}^{n}\right) \\
\left(w_{2}\right)_{j}^{n, \pm} & =\left(w_{2}\right)_{j}^{n} \pm \frac{1}{2} \operatorname{minmod}\left(\left(w_{2}\right)_{j+1}^{n}-\left(w_{2}\right)_{j}^{n},\left(w_{2}\right)_{j}^{n}-\left(w_{2}\right)_{j-1}^{n}\right)
\end{aligned}\right.
$$

which is easily seen to imply the last two constraints in (32) by definition of the minmod function (and since $\mathbf{u}_{j-1}^{n}$ and $\mathbf{u}_{j+1}^{n}$ are also assumed to be in $\Omega_{c}$ ). This shows that the "guess" choice (30) has not to be modified in order to obtain the required $L^{1}$ stability property, provided that Riemann coordinates are used in the linear reconstruction procedure.

## Accuracy in time

To conclude this section, we aim at proposing a time discretization which is second-order accurate in smooth regions and away from phase transitions. Of course, the definition of the strategy must take into account the presence of phase transitions. Actually, our objective is to propose a simple numerical time integration which is equivalent, away from phase transitions, to the well-known RK2 method (2nd order Runge-Kutta, or Heun). The latter is assumed to be familiar to the author.
The MUSCL scheme obtained below writes

$$
\overline{\mathbf{u}}_{j}^{n+1}-\mathbf{u}_{j}^{n}=\frac{\Delta x-\overline{\Delta x}_{j}^{n}}{\overline{\Delta x}_{j}^{n}} \mathbf{u}_{j}^{n}-\frac{\Delta t}{\overline{\Delta x}}{ }_{j}^{n}\left(\overline{\mathbf{f}}^{-}\left(\mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)-\overline{\mathbf{f}}^{+}\left(\mathbf{u}_{j-1}^{n,+}, \mathbf{u}_{j}^{n,-}\right)\right) .
$$

¿From this formula and the approximated values $\left(\mathbf{u}_{j}^{n}\right)_{j \in \mathbb{Z}}$ on the cells $\mathcal{C}_{j}$, we then define a first approximation $\overline{\mathbf{u}}_{j}^{n+1=}$ of the updated value on the cell $\overline{\mathcal{C}}_{j}^{n}$ by

$$
\overline{\mathbf{u}}_{j}^{n+1=}-\mathbf{u}_{j}^{n}=\frac{\Delta x-\overline{\Delta x}_{j}^{n}}{\overline{\Delta x}_{j}^{n}} \mathbf{u}_{j}^{n}-\frac{\Delta t}{\overline{\Delta x}_{j}^{n}}\left(\overline{\mathbf{f}}^{-}\left(\mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)-\overline{\mathbf{f}}^{+}\left(\mathbf{u}_{j-1}^{n,+}, \mathbf{u}_{j}^{n,-}\right)\right)
$$

Then, another one denoted $\overline{\mathbf{u}}_{j}^{n+1-}$ is obtained from the first approximations $\left(\overline{\mathbf{u}}_{j}^{n+1=}\right)_{j \in \mathbb{Z}}$ :

$$
\overline{\mathbf{u}}_{j}^{n+1-}-\mathbf{u}_{j}^{n+1=}=\frac{\Delta x-\overline{\Delta x}_{j}^{n}}{\overline{\Delta x}_{j}^{n}} \overline{\mathbf{u}}_{j}^{n+1=}-\frac{\Delta t}{\overline{\Delta x_{j}^{n}}}\left(\overline{\mathbf{f}}^{-}\left(\mathbf{u}_{j}^{n+1=,+}, \mathbf{u}_{j+1}^{n+1=,-}\right)-\overline{\mathbf{f}}^{+}\left(\mathbf{u}_{j-1}^{n+1=,+}, \mathbf{u}_{j}^{n+1=,-}\right)\right)
$$

Finally, $\overline{\mathbf{u}}_{j}^{n+1}$ is defined from these two approximations setting

$$
\overline{\mathbf{u}}_{j}^{n+1}=\mathbf{u}_{j}^{n}+\frac{1}{2}\left[\left(\overline{\mathbf{u}}_{j}^{n+1=}-\mathbf{u}_{j}^{n}\right)+\left(\overline{\mathbf{u}}_{j}^{n+1-}-\overline{\mathbf{u}}_{j}^{n+1=}\right)\right]
$$

which equivalently recasts as

$$
\begin{aligned}
\overline{\mathbf{u}}_{j}^{n+1}=\mathbf{u}_{j}^{n} & +\frac{\Delta x-\overline{\Delta x}_{j}^{n}}{2 \overline{\Delta x}_{j}^{n}}\left(\mathbf{u}_{j}^{n}+\overline{\mathbf{u}}_{j}^{n+1=}\right) \\
& -\frac{\Delta t}{2 \overline{\Delta x}_{j}^{n}}\left(\overline{\mathbf{f}}^{-}\left(\mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)-\overline{\mathbf{f}}^{+}\left(\mathbf{u}_{j-1}^{n,+}, \mathbf{u}_{j}^{n,-}\right)\right) \\
& -\frac{\Delta t}{2 \overline{\Delta x}_{j}^{n}}\left(\overline{\mathbf{f}}^{-}\left(\mathbf{u}_{j}^{n+1=,+}, \mathbf{u}_{j+1}^{n+1=,-}\right)-\overline{\mathbf{f}}^{+}\left(\mathbf{u}_{j-1}^{n+1=,+}, \mathbf{u}_{j}^{n+1=,-}\right)\right)
\end{aligned}
$$

To conclude this section, note that away from phase transitions, we have $\overline{\Delta x}_{j}^{n}=\Delta x$ and that the numerical fluxes coincide with the ones of the usual Godunov method. So that we get in such a situation

$$
\begin{aligned}
\overline{\mathbf{u}}_{j}^{n+1}=\mathbf{u}_{j}^{n} & -\frac{\Delta t}{2 \Delta x}\left(\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0 ; \mathbf{u}_{j}^{n,+}, \mathbf{u}_{j+1}^{n,-}\right)\right)-\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0 ; \mathbf{u}_{j-1}^{n,+}, \mathbf{u}_{j}^{n,-}\right)\right)\right) \\
& -\frac{\Delta t}{2 \Delta x}\left(\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0 ; \mathbf{u}_{j}^{n+1=,+}, \mathbf{u}_{j+1}^{n+1=,-}\right)\right)-\mathbf{f}\left(\mathbf{v}_{\mathbf{r}}\left(0 ; \mathbf{u}_{j-1}^{n+1=,+}, \mathbf{u}_{j}^{n+1=,-}\right)\right)\right)
\end{aligned}
$$

and the classical method consisting in a RK2 time integration together with a MUSCL reconstruction strategy for the space discretization is recovered. The scheme is then second-order accurate in both space and time in smooth regions.

## 5 Numerical experiments

We now test our algorithm on several Riemann initial data leading to typical solutions of interest, involving phase transitions or not (see also Section 3 above). The various constants involved in the models (free or congested) are taken as follows :

$$
R=1, \quad V=2, \quad V_{f}=1, \quad V_{c}=0.85, \quad Q=0.5, \quad Q^{-}=0.25, \quad Q^{+}=1.5 .
$$

Note that a state in the free phase will be characterized by the value of its density $\rho$, while for a congested state, we will use the density $\rho$ and the flux $f=\rho v_{c}(\rho, q)$ (the value of $q$ is recovered by inverting (5)). The exact and numerical profiles of the density and the speed are systematically plotted, at a final time $T_{f}$ and for a space step $\Delta x$ that will be precised for each test case.

### 5.1 The case of a single phase

In this section, we consider four Riemann initial data located in a single phase (free or congested). More precisely, the left and the right states are as follows.

Test A
$\mathbf{u}^{l}$ (free) : $\quad \rho^{l}=0.1$
$\mathbf{u}^{r}$ (free) : $\quad \rho^{r}=0.4$

Test B

$$
\begin{array}{lr}
\hline \mathbf{u}^{l} \text { (free) : } & \rho^{l}=0.4 \\
\mathbf{u}^{r} \text { (free) }: & \rho^{r}=0.25
\end{array}
$$

Test C

| $\mathbf{u}^{l}$ (congested) : | $\rho^{l}=0.7$ | $f^{l}=0.2$ |
| :---: | :---: | :---: |
| $\mathbf{u}^{r}$ (congested) : | $\rho^{r}=0.4$ | $f^{r}=0.3$ |

Test D

$$
\begin{array}{ccc}
\hline \mathbf{u}^{l} \text { (congested) : } & \rho^{l}=0.4 & f^{l}=0.3 \\
\mathbf{u}^{r} \text { (congested) : } & \rho^{r}=0.7 & f^{r}=0.2
\end{array}
$$

In all these situations, there is no phase transition. Then, our method simply reduces to the usual Godunov method since at each interface $x_{j+1 / 2}, j \in \mathbb{Z}$, we necessarily have $\sigma_{j+1 / 2}^{n}=0$. In particular, the method is conservative. In the free phase, the left state is connected directly to the right state, by a shock in Test A (see Figure 9) and a rarefaction in Test B (see Figure 10). In the congested phase, the solution is generally made of two distinct waves. More precisely, we have a rarefaction wave followed by a contact discontinuity in the case of Test C (Figure 11), and a shock wave followed by a contact discontinuity in the case of Test D (Figure 12). Solutions are plotted at time $T_{f}=0.4$ for Test A, C and $T_{f}=0.5$ for Tests B, D. The computations have been performed with a mesh containing 100 points per unit interval ( $\Delta x=0.01$ ). As expected, we observe that the numerical solutions are in good agreement with exact ones.


Figure 9: Test A: $\rho$ (Left) and $v$ (Right)


Figure 10: Test B : $\rho$ (Left) and $v$ (Right)


Figure 11: Test C : $\rho$ (Left) and $v$ (Right)


Figure 12: Test D : $\rho$ (Left) and $v$ (Right)

### 5.2 The case of a phase transition from congested to free

In order to assess the validity of our method, we now consider various situations where a phase transition is present in the solution. We begin with phase transitions from a congested state to a free state.

For the Test $\mathbf{E}$, we consider the following left and right states for the Riemann initial data

Test E
$\mathbf{u}^{l}$ (congested) : $\quad \rho^{l}=0.7 \quad f^{l}=0.3$
$\mathbf{u}^{r}$ (free) : $\quad \rho^{r}=0.3$
leading to a solution made of a rarefaction in the congested phase, followed by a phase transition to a free state, itself followed by a rarefaction wave in the free phase (see Figure 3, left). The solutions are plotted on Figure 13 at time $T_{f}=0.5$. For this test case, we have used a mesh containing 500 points ( $\Delta x=0.002$ ).

The next table concerns Test F.

Test F
$\mathbf{u}^{l}$ (congested) : $\quad \rho^{l}=0.45 \quad f^{l}=0.25$
$\mathbf{u}^{r}$ (free) : $\quad \rho^{r}=0.3$

Here, the solution is a shock-like phase transition from the left state to a free state, followed by a shock wave in the free phase (see Figure 3, right). See Figure 14 for the solution at time $T_{f}=0.35$ with $\Delta x=0.01$.


Figure 13: Test E : $\rho$ (Left) and $v$ (Right)


Figure 14: Test F : $\rho$ (Left) and $v$ (Right)

For these two test cases, we observe a very good agreement between the exact and numerical solutions. In particular, we note that the phase transitions we obtain join the right states and propagate apparently with the right speed (theoretically given by Rankine-Hugoniot conditions). Which means that the mass conservation property is not lost at the discrete level.
In the case of Test $\mathbf{E}$, we note that the numerical diffusion of the scheme on the rarefaction waves makes the free state of the phase transition hardly apparent at the final time $T_{f}$ proposed. This phenomenon will be even more pronounced in the numerical solutions presented below for Tests H, I.

### 5.3 The case of a phase transition from free to congested

Let us now address the case of phase transitions from a free state to a congested state.
For the Test G, we choose

Test G

| $\mathbf{u}^{l}$ (free) : | $\rho^{l}=0.35$ |
| ---: | :---: |
| $\mathbf{u}^{r}($ congested $):$ | $\rho^{r}=0.6 \quad f^{r}=0.25$ |

The corresponding solution is a shock-like phase transition to a congested state followed by a contact discontinuity (see Figure 4 , left). Figure 15 plots the solution at time $T_{f}=0.6$ with $\Delta x=0.002$.


Figure 15: Test G: $\rho$ (Left) and $v$ (Right)

For Test H and Test I, we take

Test H

$$
\begin{array}{cc}
\mathbf{u}^{l} \text { (free) }: & \rho^{l}=0.24 \\
\mathbf{u}^{r}(\text { congested }): & \rho^{r}=0.7 \quad f^{r}=0.2
\end{array}
$$

Test I

$$
\begin{array}{rc}
\mathbf{u}^{l} \text { (free) }: & \rho^{l}=0.215 \\
\mathbf{u}^{r}(\text { congested }): & \rho^{r}=0.7 \quad f^{r}=0.2
\end{array}
$$

leading to two solutions composed of a phase transition followed by a rarefaction wave, and a contact discontinuity propagating with a positive speed. In the case of Test H, the "foot" of the rarefaction wave propagates with a speed very close to the one of the phase transition (see Figure 4, right), while for Test I, the rarefaction is simply attached to the transition (see Figure 5, left). These properties make the congested state of the phase transition more difficult to capture properly, due to the numerical diffusion of the scheme which is present in the rarefaction wave. Note that this state is always over-estimated from the proposed averaging strategy. However we observe a good agreement between the numerical solution and the exact solution, and as it is expected, the numerical solution becomes really better when the order of accuracy of the method is higher as it is illustrated on Figures 16 and 17 where $\Delta x=0.0005$ ( $T_{f}=0.8$ for both cases).


Figure 16: Test H: $\rho$ (Left) and $v$ (Right)


Figure 17: Test I : $\rho$ (Left) and $v$ (Right)

At last, we consider the Test $\mathbf{J}$ for which the solution is a phase transition followed by a contact discontinuity (see Figure 5, right). The initial data is such that

## Test J

| $\mathbf{u}^{l}$ (free) : | $\rho^{l}=0.1$ |
| :---: | :---: |
| $\mathbf{u}^{r}($ congested $):$ | $\rho^{r}=0.7 \quad f^{r}=0.2$ |

and the solutions are plotted on Figure 18 at time $T_{f}=1.5(\Delta x=0.01)$.


Figure 18: Test J : $\rho$ (Left) and $v$ (Right)

We have covered in this section all possible situations described in Section 3 for a Riemann problem. We have observed that the numerical solutions are in accordance with the exact ones and exhibit sharp phase transitions when present (i.e. the corresponding profiles are infinitely thin). This proves the relevance of our approach. In the next section, we propose to estimate the mass conservation errors introduced by the sampling procedure of the scheme. Of course, these are expected to be very small since otherwise the numerical solutions would not be as close as they are to the exact ones. Moreover, we will see that the conservation errors decrease with the mesh size.

## 6 Measure of the conservation errors

Due to the random sampling present in Step 2 (Modified), our method does not strictly conserve the mass $\rho$. We then propose to measure the conservation errors on piecewise constant numerical solution $\rho_{\nu}$ defined as

$$
\rho_{\nu}(x, t)=\rho_{j}^{n} \quad \text { if } \quad(x, t) \in\left[x_{j-1 / 2}, x_{j+1 / 2}\right) \times\left[t^{n}, t^{n+1}\right),
$$

between times $t=0$ and $t=T$, for some $T>0$. We denote $\left[x_{0}, x_{1}\right]$ the computational domain and we proceed exactly as in [4] : we compare with 0 the function $E: T \in \mathbb{R}^{+} \rightarrow E(T) \in \mathbb{R}$ with $E(T)$ defined by relation

$$
\begin{gather*}
\int_{x_{0}}^{x_{1}} \rho_{\nu}(x, T) d x \times E(T)= \\
\int_{x_{0}}^{x_{1}} \rho_{\nu}(x, T) d x-\int_{x_{0}}^{x_{1}} \rho_{\nu}(x, 0) d x+\int_{0}^{T}\left\{\rho v_{c}(\rho, q)\right\}_{\nu}\left(x_{1}, t\right) d t-\int_{0}^{T}\left\{\rho v_{c}(\rho, q)\right\}_{\nu}\left(x_{0}, t\right) d t \tag{33}
\end{gather*}
$$

Recall that $q=\rho V$ in the free phase. $E(T)$ represents the relative conservation error of $\rho$ at time $T$ on interval $\left[x_{0}, x_{1}\right]$. In the next tables, we give for the Tests $\mathbf{E}, \mathbf{F}, \mathbf{G}, \mathbf{H}, \mathbf{I}, \mathbf{J}$ the values of the $L^{1}$-norm $\frac{1}{T_{f}}\|E\|_{L^{1}\left(0, T_{f}\right)}$ of $E$, namely

$$
\frac{1}{T_{f}}\|E\|_{L^{1}\left(0, T_{f}\right)}=\frac{1}{T_{f}} \int_{0}^{T_{f}}|E(T)| d T=\sum_{t^{n}=0}^{t^{n+1}=T_{f}} \frac{\left(t^{n+1}-t^{n}\right)}{T_{f}}\left|E\left(t^{n}\right)\right|,
$$

where $T_{f}$ is the final time of the corresponding simulations presented in the previous section. We will consider four different meshes containing 100, 500, 1000 and 2000 points per unit interval. The computational domain $\left[x_{0}, x_{1}\right]$ is always $[-0.5,0.5]$ except for Test $\mathbf{J}$ for which it is $[-0.2,0.8]$. Of course, when the solution remains in a single phase (Tests A, B, C, D), our algorithm reduces to the classical Godunov scheme and so is actually conservative. As expected, we observe for the other test cases that the conservation errors are decreasing with the mesh size. We do not note a great improvement when dealing with the second-order scheme. This is not really surprising since conservation errors come from the presence of phase transitions, and the scheme remains first-order accurate near phase transitions.

| \# of points | Test E | Test F | Test G | Test H | Test I | Test J |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 100 | $0.44 \%$ | $0.22 \%$ | $0.64 \%$ | $0.39 \%$ | $0.91 \%$ | $0.65 \%$ |
| 500 | $0.16 \%$ | $0.11 \%$ | $0.17 \%$ | $0.11 \%$ | $0.22 \%$ | $0.15 \%$ |
| 1000 | $0.094 \%$ | $0.075 \%$ | $0.095 \%$ | $0.055 \%$ | $0.11 \%$ | $0.081 \%$ |
| 2000 | $0.051 \%$ | $0.039 \%$ | $0.057 \%$ | $0.025 \%$ | $0.052 \%$ | $0.045 \%$ |

Figure 19: Conservation errors (first-order scheme).

| \# of points | Test E | Test F | Test G | Test H | Test I | Test J |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 100 | $0.25 \%$ | $0.26 \%$ | $0.23 \%$ | $0.35 \%$ | $0.87 \%$ | $0.71 \%$ |
| 500 | $0.054 \%$ | $0.12 \%$ | $0.071 \%$ | $0.10 \%$ | $0.21 \%$ | $0.19 \%$ |
| 1000 | $0.030 \%$ | $0.08 \%$ | $0.044 \%$ | $0.054 \%$ | $0.11 \%$ | $0.11 \%$ |
| 2000 | $0.016 \%$ | $0.041 \%$ | $0.031 \%$ | $0.027 \%$ | $0.052 \%$ | $0.05 \%$ |

Figure 20: Conservation errors (second-order scheme).
For the sake of completeness, we give on the following two tables the $L^{1}$ error between the numerical solution and the exact one for all the test cases and several mesh sizes. Note that the $L^{1}$ norm is computed on the whole computational domain, that is taking into account regions that are possibly not smooth. However, we observe that for each scheme the errors decrease with the mesh size (the convergence seems secured), while the second-order scheme allows to obtain smaller errors as it is expected.

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| \# of points | Test A | Test B | Test C | Test D | Test E |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 100 | $2.29 e^{-3}$ | $3.22 e^{-3}$ | $7.87 e^{-3}$ | $9.50 e^{-3}$ | $8.64 e^{-3}$ |
| 500 | $4.58 e^{-4}$ | $9.87 e^{-4}$ | $3.17 e^{-3}$ | $4.29 e^{-3}$ | $2.99 e^{-3}$ |
| 1000 | $2.29 e^{-4}$ | $5.72 e^{-4}$ | $2.08 e^{-3}$ | $3.04 e^{-3}$ | $1.74 e^{-3}$ |
| 2000 | $1.14 e^{-4}$ | $3.26 e^{-4}$ | $1.34 e^{-3}$ | $2.15 e^{-3}$ | $1.05 e^{-3}$ |
| \# of points | Test F | Test G | Test H | Test I | Test J |
| 100 | $3.50 e^{-3}$ | $9.67 e^{-3}$ | $9.84 e^{-3}$ | $8.19 e^{-3}$ | $1.18 e^{-2}$ |
| 500 | $7.76 e^{-4}$ | $2.03 e^{-3}$ | $4.81^{-3}$ | $4.76 e^{-3}$ | $5.13 e^{-3}$ |
| 1000 | $6.05 e^{-4}$ | $1.63 e^{-3}$ | $3.55 e^{-3}$ | $3.27 e^{-3}$ | $2.82 e^{-3}$ |
| 2000 | $1.95 e^{-4}$ | $1.01 e^{-3}$ | $2.53 e^{-3}$ | $2.12 e^{-3}$ | $2.13 e^{-3}$ |

Figure 21: $L^{1}$ errors (first-order scheme).

| \# of points | Test A | Test B | Test C | Test D | Test E |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 100 | $1.73 e^{-3}$ | $1.07 e^{-3}$ | $4.27 e^{-3}$ | $5.92 e^{-3}$ | $4.18 e^{-3}$ |
| 500 | $3.48 e^{-4}$ | $2.16 e^{-4}$ | $1.11 e^{-3}$ | $2.11 e^{-3}$ | $8.60 e^{-4}$ |
| 1000 | $1.74 e^{-4}$ | $1.08 e^{-4}$ | $6.10 e^{-4}$ | $1.34 e^{-3}$ | $3.17 e^{-4}$ |
| 2000 | $8.69 e^{-5}$ | $5.40 e^{-5}$ | $3.37 e^{-4}$ | $8.52 e^{-4}$ | $2.15 e^{-4}$ |
| \# of points | Test F | Test G | Test H | Test I | Test J |
| 100 | $3.00 e^{-3}$ | $5.12 e^{-3}$ | $7.42 e^{-3}$ | $5.00 e^{-3}$ | $7.98 e^{-3}$ |
| 500 | $5.98 e^{-4}$ | $6.19 e^{-4}$ | $2.16 e^{-3}$ | $2.48 e^{-3}$ | $3.02 e^{-3}$ |
| 1000 | $5.16 e^{-4}$ | $7.23 e^{-4}$ | $1.25 e^{-3}$ | $1.41 e^{-3}$ | $1.23 e^{-3}$ |
| 2000 | $1.49 e^{-4}$ | $4.13 e^{-4}$ | $7.59 e^{-4}$ | $6.97 e^{-4}$ | $7.12 e^{-4}$ |

Figure 22: $L^{1}$ errors (second-order scheme).

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