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# Godunov type scheme for the linear wave equation with Coriolis source term 

Emmanuel Audusse, Stéphane Dellacherie, Do Minh Hieu, Pascal Omnes, Yohan Penel

May 31, 2016


#### Abstract

We propose a method to explain the behaviour of the Godunov finite volume scheme applied to the linear wave equation with Coriolis source term at low Froude number. In particular, we use the Hodge decomposition and we study the properties of the modified equation associated to the Godunov scheme. Based on the structure of the discrete kernel of the linear operator discretized by using the Godunov scheme, we clearly explain the inaccuracy of the classical Godunov scheme at low Froude number and we introduce a way to modify it to recover a correct accuracy.


## 1 Introduction

In this communication, we study some procedures to make finite volume Godunov type schemes accurate when solving perturbations around a steady-state. In what follows, we restrict the analysis to the quasi-1d linear wave equation with Coriolis term. Nevertheless, our future main objective is to derive accurate and stable finite volume collocated schemes for the dimensionless shallow water equations

$$
\left\{\begin{array}{l}
\text { St } \partial_{t} h+\nabla \cdot(h \overline{\mathbf{u}})=0,  \tag{1a}\\
\text { St } \partial_{t}(h \overline{\mathbf{u}})+\nabla \cdot(h \overline{\mathbf{u}} \otimes \overline{\mathbf{u}})+\frac{1}{\mathrm{Fr}^{2}} \nabla\left(\frac{h^{2}}{2}\right)=-\frac{1}{\mathrm{Fr}^{2}} h \nabla b-\frac{1}{\mathrm{Ro}} h \overline{\mathbf{u}}^{\perp},
\end{array}\right.
$$

in a rotating frame when the flow is a perturbation around the so-called geostrophic equilibrium. In System (1) unknowns $h$ and $\overline{\mathbf{u}}$ respectively denote the water depth and the velocity of the water column and function $b(x)$ denotes the topography of the considered oceanic basin and is a given function. Dimensionless numbers St, Fr and Ro respectively stand for the Strouhal, the Froude and the Rossby numbers defined by

$$
\text { St }=\frac{L}{U T}, \quad \operatorname{Fr}=\frac{U}{\sqrt{g H}}, \quad \mathrm{Ro}=\frac{U}{\Omega L}
$$

where the parameter $g$ and $\Omega$ denote the gravity coefficient and the angular velocity of the Earth. Constants $U, H, L$ and $T$ are some characteristic velocity, vertical and horizontal lengths and time. In the sequel, we shall focus on cases where

$$
\mathrm{Ro}=\mathcal{O}(M) \quad \text { and } \quad \mathrm{Fr}=\mathcal{O}(M)
$$

with $M$ a small parameter. For large scale oceanographic flows, typical values lead to $M \sim 10^{-2}$ since

$$
U \approx 1 \mathrm{~m} \cdot \mathrm{~s}^{-1}, \quad L \approx 10^{6} \mathrm{~m}, \quad H \approx 10^{3} \mathrm{~m}, \quad \Omega \approx 10^{-4} \mathrm{rad} \cdot \mathrm{~s}^{-1}
$$

At the leading order, the stationary state associated to this equation (or let us say the solution for Strouhal number of order 1) satisfies the so-called lake at rest equilibrium

$$
\begin{equation*}
\nabla(h+b)=0 \tag{2}
\end{equation*}
$$

At the next order, the flow satisfies the so-called geostrophic equilibrium

$$
\begin{equation*}
\nabla h=-\overline{\mathbf{u}}^{\perp} . \tag{3}
\end{equation*}
$$

Note that this relation implies

$$
\begin{equation*}
\nabla \cdot \overline{\mathbf{u}}=0 \tag{4}
\end{equation*}
$$

The ability of numerical schemes to well capture the particular solutions (2) and (3) is of great practical interest since it has a direct consequence on the accuracy of the numerical solution when perturbations around these equilibria are considered. A huge litterature has been devoted to the preservation of the lake at rest equilibrium (2), see in particular [3] and references therein. The question of the geostrophic equilibrium (3) including the divergence constraint (4) is more complex. It has been studied in a finite element framework by Le Roux [11]. The author considers in his work the linearized version of System (1) and studies the behaviour of several types of finite elements. He shows that spurious modes are created, in particular when the number of degrees of freedom is not the same for height and velocity. In the finite volume framework, the non linear case has been studied in [4,6,13]. In particular, Bouchut and coauthors introduce in [4] the apparent topography method that allows to adapt to this problem the hydrostatic reconstruction method [2] that was developed to ensure the preservation of the lake at rest equilibrium (2).

Let us now suppose that the topography is flat and the solution is independent of the $y$ direction. For a Strouhal number of order $\mathcal{O}\left(M^{-1}\right)$, i.e. for short times and for Rossby and Froude numbers of order $\mathcal{O}(M)$, the solution of system (1) satisfies at the leading order the quasi-1d linear wave equation with Coriolis term

$$
\left\{\begin{array}{l}
\partial_{t} r+a_{\star} \partial_{x} u=0  \tag{5}\\
\partial_{t} u+a_{\star} \partial_{x} r=\omega v \\
\partial_{t} v=-\omega u
\end{array}\right.
$$

where $a_{\star}$ and $\omega$ are constants of order one, respectively related to the wave velocity and to the rotating velocity. The stationary state corresponding to Equation (5) is the 1d version of the geostrophic equilibrium (3) and is called $1 D$ geostrophic equilibrium. It is such that

$$
\begin{equation*}
u=0, \quad a_{\star} \partial_{x} r=\omega v . \tag{6}
\end{equation*}
$$

Many works were devoted to the study of the homogeneous wave equations. In particular, in a serie of articles [7-10], Dellacherie and coauthors studied the behaviour of Godunov type schemes for the 2 d linear wave equation. Their works are part of a more general study about the use of Godunov type schemes in the context of the incompressible limit for Euler equations, i.e. for low Mach number flows, see for example [15-17]. In the present work, we extend the aforementioned approach from Dellacherie and coauthors to take into account the Coriolis source term. First, in Section 2, we analyze the continuous case by using a Hodge type decomposition. Then, in Section 3 and 4, we study three Godunov type numerical schemes to compute approximate solutions of Equation (5):

- The Classical Godunov scheme;
- The Low Froude Godunov scheme;
- The All Froude Godunov scheme.

For each scheme, we study the kernel of the discrete operator and we compare it to the continuous kernel (6). Then, we study the accuracy of the scheme at low Froude number, i.e. when the initial solution is close to the kernel. This is done first for the modified equation associated to the scheme in Section 3 and then in the fully discrete case in Section 4. Moreover we study the stability of each discrete scheme by using Fourier analysis. Finally we present in Section 5 some numerical results to illustrate our purpose.

## 2 Properties of the linear wave equation with Coriolis source term

We first focus on the properties of the linear wave equation on the 1 d torus $\mathbb{T}$. To begin with, we introduce the Hilbert space

$$
\left(L^{2}(\mathbb{T})\right)^{3}=\left\{q=(r, u, v) \mid \int_{\mathbb{T}} r^{2} \mathrm{~d} x+\int_{\mathbb{T}}\left(u^{2}+v^{2}\right) \mathrm{d} x<\infty\right\}
$$

equipped with the scalar product

$$
\left\langle q_{1}, q_{2}\right\rangle=\int_{\mathbb{T}} r_{1} r_{2} \mathrm{~d} x+\int_{\mathbb{T}}\left(u_{1} u_{2}+v_{1} v_{2}\right) \mathrm{d} x
$$

### 2.1 Structure of the kernel of the original model

Let us define the following space

$$
\begin{equation*}
\mathcal{E}_{\omega \neq 0}=\left\{q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3} \mid u=0, \forall \phi \in C_{c}^{\infty}(\mathbb{T}), \int_{\mathbb{T}} a_{\star} r \partial_{x} \phi \mathrm{~d} x=-\int_{\mathbb{T}} \omega v \phi \mathrm{~d} x\right\} \tag{7}
\end{equation*}
$$

We then prove this preliminary result:
Lemma 1. The orthogonal of $\mathcal{E}_{\omega \neq 0}$ is

$$
\mathcal{E}_{\omega \neq 0}^{\perp}=\left\{q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3} \mid \forall \varphi \in C_{c}^{\infty}(\mathbb{T}), \int_{\mathbb{T}} a_{\star} v \partial_{x} \varphi \mathrm{~d} x=-\int_{\mathbb{T}} \omega r \varphi \mathrm{~d} x\right\} .
$$

Moreover, we have $\mathcal{E}_{\omega \neq 0} \oplus \mathcal{E}_{\omega \neq 0}^{\perp}=\left(L^{2}(\mathbb{T})\right)^{3}$. In other words, any $q \in\left(L^{2}(\mathbb{T})\right)^{3}$ can be uniquely decomposed into

$$
\begin{equation*}
q=\hat{q}+\tilde{q} \tag{8}
\end{equation*}
$$

where $\hat{q} \in \mathcal{E}_{\omega \neq 0}$ and $\tilde{q} \in \mathcal{E}_{\omega \neq 0}^{\perp}$.
The Hodge decomposition (8) allows us to define the orthogonal projection

$$
\mathbb{P}:\left\{\begin{align*}
\left(L^{2}(\mathbb{T})\right)^{3} & \longrightarrow \mathcal{E}_{\omega \neq 0}  \tag{9}\\
q & \longmapsto \hat{q}
\end{align*}\right.
$$

Remark 1. The kernel and its orthogonal set can be described in a simpler way due to the definition of Sobolev spaces, namely

$$
\begin{aligned}
& \mathcal{E}_{\omega \neq 0}=\left\{q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3} \mid r \in H^{1}(\mathbb{T}), u=0, v=\frac{a_{\star}}{\omega} r^{\prime}\right\}, \\
& \mathcal{E}_{\omega \neq 0}^{\perp}=\left\{q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3} \mid v \in H^{1}(\mathbb{T}), r=\frac{a_{\star}}{\omega} v^{\prime}\right\} .
\end{aligned}
$$

Moreover, the fact that we consider periodic functions implies that for $\hat{q} \in \mathcal{E}_{\omega \neq 0}$ and $\tilde{q} \in \mathcal{E}_{\omega \neq 0}^{\perp}$, we have

$$
\int_{\mathbb{T}} \hat{v} \mathrm{~d} x=0 \quad \text { and } \quad \int_{\mathbb{T}} \tilde{r} \mathrm{~d} x=0
$$

due to boundary conditions.
Proof. Our purpose is to prove that $\mathcal{E}_{\omega \neq 0}^{\perp}=A$ where

$$
A=\left\{\left.q=\left(\frac{a_{\star}}{\omega} v^{\prime}, u, v\right) \right\rvert\, u \in L^{2}(\mathbb{T}), v \in H^{1}(\mathbb{T})\right\} .
$$

Firstly, let us prove that $A \subset \mathcal{E}_{\omega \neq 0}^{\perp}$. For $\tilde{q} \in A$, we have

$$
\forall q \in \mathcal{E}_{\omega \neq 0},\langle\tilde{q}, q\rangle=\int_{\mathbb{T}} r \frac{a_{\star}}{\omega} \tilde{v}^{\prime} \mathrm{d} x+\int_{\mathbb{T}} \frac{a_{\star}}{\omega} r^{\prime} \tilde{v} \mathrm{~d} x=\frac{a_{\star}}{\omega}\left(\int_{\mathbb{T}} r \tilde{v}^{\prime} \mathrm{d} x+\int_{\mathbb{T}} r^{\prime} \tilde{v} \mathrm{~d} x\right) .
$$

According to [5, Corollary 8.10] with $(r, \tilde{v}) \in\left(H^{1}(\mathbb{T})\right)^{2}$, we have $r \tilde{v} \in H^{1}(\mathbb{T})$ and $(r \tilde{v})^{\prime}=r^{\prime} \tilde{v}+r \tilde{v}^{\prime}$. Therefore, we obtain, thanks to periodic boundary conditions on $\mathbb{T}$

$$
\int_{\mathbb{T}} r \tilde{v}^{\prime} \mathrm{d} x+\int_{\mathbb{T}} r^{\prime} \tilde{v} \mathrm{~d} x=\int_{\mathbb{T}}(r \tilde{v})^{\prime} \mathrm{d} x=0
$$

which leads to $\langle\tilde{q}, q\rangle=0, \forall q \in \mathcal{E}_{\omega \neq 0}$. It means that $\tilde{q} \in \mathcal{E}_{\omega \neq 0}^{\perp}$.
Secondly, we prove that $\mathcal{E}_{\omega \neq 0}^{\perp} \subset A$. Let $\tilde{q} \in \mathcal{E}_{\omega \neq 0}^{\perp}$. Therefore

$$
\forall r \in H^{1}(\mathbb{T}), \int_{\mathbb{T}} \tilde{r} r \mathrm{~d} x+\int_{\mathbb{T}} \frac{a_{\star}}{\omega} \tilde{v} r^{\prime} \mathrm{d} x=0,
$$

which implies

$$
\forall r \in C_{c}^{\infty}(\mathbb{T}), \int_{\mathbb{T}} \tilde{v} r^{\prime} \mathrm{d} x=-\frac{\omega}{a_{\star}} \int_{\mathbb{T}} \tilde{r} r \mathrm{~d} x
$$

As a result, $\tilde{v} \in H^{1}(\mathbb{T})$ and $a_{\star} \tilde{v}^{\prime}=\omega \tilde{r}$. We come to the conclusion that

$$
\mathcal{E}_{\omega \neq 0}^{\perp}=A=\left\{q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3} \mid \forall \varphi \in C_{c}^{\infty}(\mathbb{T}), \int_{\mathbb{T}} a_{\star} v \partial_{x} \varphi \mathrm{~d} x=-\int_{\mathbb{T}} \omega r \varphi \mathrm{~d} x\right\} .
$$

We eventually have to prove that

$$
\mathcal{E}_{\omega \neq 0} \oplus \mathcal{E}_{\omega \neq 0}^{\perp}=\left(L^{2}(\mathbb{T})\right)^{3} .
$$

We only have to check $\left(L^{2}(\mathbb{T})\right)^{3} \subset \mathcal{E}_{\omega \neq 0} \oplus \mathcal{E}_{\omega \neq 0}^{\perp}$, because of the fact that $\mathcal{E}_{\omega \neq 0} \oplus \mathcal{E}_{\omega \neq 0}^{\perp} \subset\left(L^{2}(\mathbb{T})\right)^{3}$ is trivial.
For $q=(r, u, v) \in\left(L^{2}(\mathbb{T})\right)^{3}$, let us set

$$
\begin{aligned}
\hat{r} & =\mu(r)-h, & \tilde{r} & =r-\mu(r)+h, \\
\hat{u} & =0, & \tilde{u} & =u, \\
\hat{v} & =-\frac{a_{\star}}{\omega} \partial_{x} h, & \partial_{x} \tilde{v} & =\frac{\omega}{a_{\star}}(r-\mu(r)+h) \text { and } \int_{\mathbb{T}} \tilde{v} \mathrm{~d} x=\int_{\mathbb{T}} v \mathrm{~d} x,
\end{aligned}
$$

where $\mu(r)=\frac{1}{|\mathbb{T}|} \int_{\mathbb{T}} r \mathrm{~d} x$ and $h \in H^{1}(\mathbb{T})$ is the unique solution of the variational formulation

$$
\forall \varphi \in H^{1}(\mathbb{T}), \int_{T} \partial_{x} h \partial_{x} \varphi \mathrm{~d} x+\frac{\omega^{2}}{a_{\star}^{2}} \int_{\mathbb{T}} \varphi h \mathrm{~d} x=-\frac{\omega}{a_{\star}} \int_{\mathbb{T}} v \partial_{x} \varphi \mathrm{~d} x-\frac{\omega^{2}}{a_{\star}^{2}} \int_{\mathbb{T}}(r-\mu(r)) \varphi \mathrm{d} x .
$$

The existence and uniqueness of $h \in H^{1}(\mathbb{T})$ results from the Lax-Milgram theorem for $\omega \neq 0$.
We easily check that $\hat{q} \in \mathcal{E}_{\omega \neq 0}$ and $\tilde{q} \in \mathcal{E}_{\omega \neq 0}^{\perp}$. To reach the conclusion, we have to check that $\hat{q}+\tilde{q}=q$. The equalities $\hat{r}+\tilde{r}=r$ and $\hat{u}+\tilde{u}=u$ are trivially verified. For $v$, we have:

$$
\forall \Phi \in C^{\infty}(\mathbb{T}), \int_{\mathbb{T}}(v-\hat{v}-\tilde{v}) \partial_{x} \Phi \mathrm{~d} x=\int_{\mathbb{T}} v \partial_{x} \Phi \mathrm{~d} x+\frac{a_{\star}}{\omega} \int_{\mathbb{T}} \partial_{x} h \partial_{x} \Phi \mathrm{~d} x+\frac{\omega}{a_{\star}} \int_{\mathbb{T}}(r-\mu(r)+h) \Phi \mathrm{d} x=0
$$

due to the choice of $h$. Using the density of $C^{\infty}(\mathbb{T})$ in $H^{1}(\mathbb{T})$, we obtain that $v-(\hat{v}+\tilde{v})=c$. By using the fact that $\int_{\mathbb{T}} \hat{v} d x=0$ and $\int_{\mathbb{T}} \tilde{v} d x=\int_{\mathbb{T}} v d x$, we get $c=0$. Therefore, we have $\hat{v}+\tilde{v}=v$.

### 2.2 Behaviour of the solution

By using Lemma 1, we obtain the following properties for the linear wave equation (5):
Proposition 1. Let $q$ be a solution of (5) on $\mathbb{T}$ with initial condition $q^{0}$. Then:
i. $\forall q^{0} \in \mathcal{E}_{\omega \neq 0}$, we have $q(t>0, \cdot)=q^{0} \in \mathcal{E}_{\omega \neq 0}$.
ii. $\forall q^{0} \in \mathcal{E}_{\omega \neq 0}^{\perp}$, we have $q(t>0, \cdot) \in \mathcal{E}_{\omega \neq 0}^{\perp}$.

Proof. We note that System (5) can be written as

$$
\partial_{t} q+A \partial_{x} q+B q=0 \quad \text { where } \quad A=\left(\begin{array}{ccc}
0 & a_{\star} & 0 \\
a_{\star} & 0 & 0 \\
0 & 0 & 0
\end{array}\right) \quad \text { and } \quad B=\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & 0 & -\omega \\
0 & \omega & 0
\end{array}\right)
$$

Due to the fact that matrix A has 3 real distinct eigenvalues $\left(\lambda=0, \lambda=-a_{\star}\right.$ and $\left.\lambda=a_{\star}\right)$, System (5) is strictly hyperbolic. Therefore this system has a unique solution [1, Th. 2.22]. And for any initial condition $q^{0}=\left(r^{0}, u^{0}, v^{0}\right)$ in $\mathcal{E}_{\omega \neq 0}$, it is obvious that this unique solution is given by $q(t>0, \cdot)=q^{0}$, which proves (i).

Let $q^{0}=\left(r^{0}, u^{0}, v^{0}\right) \in \mathcal{E}_{\omega \neq 0}^{\perp}$. We notice that

$$
\left\{\begin{array}{l}
r=r^{0}-a_{\star} \int_{0}^{t} \partial_{x} u \mathrm{~d} \tau \\
u=u^{0}-\int_{0}^{t}\left(a_{\star} \partial_{x} r-\omega v\right) \mathrm{d} \tau \\
v=v^{0}-\int_{0}^{t} \omega u \mathrm{~d} \tau
\end{array}\right.
$$

Therefore, for all $\hat{q} \in \mathcal{E}_{\omega \neq 0}$, we obtain

$$
\begin{aligned}
\langle q, \hat{q}\rangle & =\left\langle q^{0}, \hat{q}\right\rangle-a_{\star} \int_{\mathbb{T}} \int_{0}^{t} \partial_{x} u \hat{r} \mathrm{~d} x \mathrm{~d} \tau-\int_{\mathbb{T}} \int_{0}^{t} \omega u \hat{v} \mathrm{~d} x \mathrm{~d} \tau=\left\langle q^{0}, \hat{q}\right\rangle-\int_{0}^{t} \int_{\mathbb{T}} a_{\star} \partial_{x} u \hat{r} \mathrm{~d} x \mathrm{~d} \tau-\int_{0}^{t} \int_{\mathbb{T}} \omega u \hat{v} \mathrm{~d} x \mathrm{~d} \tau \\
& =\left\langle q^{0}, \hat{q}\right\rangle+\int_{0}^{t} \int_{\mathbb{T}} a_{\star} \partial_{x} \hat{r} u \mathrm{~d} x \mathrm{~d} \tau-\int_{0}^{t} \int_{\mathbb{T}} \omega \hat{v} u \mathrm{~d} x \mathrm{~d} \tau=\left\langle q^{0}, \hat{q}\right\rangle=0
\end{aligned}
$$

As a result, we conclude that $q(t>0, \cdot) \in \mathcal{E}_{\omega \neq 0}^{\perp}$, which proves (ii).
Corollary 1. Let $q$ be the solution of (5) with initial condition $q^{0}$. Then, $q$ can be decomposed into

$$
q=\mathbb{P} q^{0}+\left(q-\mathbb{P}^{0}\right) \in \mathcal{E}_{\omega \neq 0} \oplus \mathcal{E}_{\omega \neq 0}^{\perp}
$$

### 2.3 Evolution of the energy

Let us define the energy as $E=\langle q, q\rangle$.
Proposition 2. Let $q$ be the solution of (5) on $\mathbb{T}$. Then, the energy is conserved

$$
E(t>0)=E(t=0)
$$

Proof. Because $q$ is the solution of (5), we have

$$
\left\{\begin{array}{l}
\partial_{t} r=-a_{\star} \partial_{x} u \\
\partial_{t} u=\omega v-a_{\star} \partial_{x} r \\
\partial_{t} v=-\omega u
\end{array}\right.
$$

which allows to obtain

$$
\begin{aligned}
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t}\langle q, q\rangle & =a_{\star} \int_{\mathbb{T}} r\left(-\partial_{x} u\right) \mathrm{d} x+\int_{\mathbb{T}} u\left(\omega v-a_{\star} \partial_{x} r\right) \mathrm{d} x+\int_{\mathbb{T}} v(-\omega u) \mathrm{d} x \\
& =a_{\star} \int_{\mathbb{T}} r\left(-\partial_{x} u\right) \mathrm{d} x+a_{\star} \int_{\mathbb{T}} u\left(-\partial_{x} r\right) \mathrm{d} x=0
\end{aligned}
$$

Hence we have $E^{\prime}(t)=0$ which concludes the proof.
Corollary 2. For all times $t>0$, we have $\left\|q(t, \cdot)-\mathbb{P} q^{0}\right\|=\left\|q^{0}-\mathbb{P} q^{0}\right\|$.

## 3 Properties of the first order modified equation associated to the Godunov finite volume scheme

It is well known that the classical Godunov scheme is inaccurate at low Mach number (or low Froude number). With the homogeneous linear wave equation $(\omega=0)$, the problem appears only in the 2 d case over rectangular meshes. The work in $[8,10]$ clearly points out the main reason of the inaccuracy. Shortly, this is because the classical Godunov scheme suffers from the loss of invariance of the well-prepared subspace $\mathcal{E}$ when the numerical diffusion related to the velocity equation is not equal to 0 . However in our case with Coriolis source term, the problem appears already in $1 d$ due to the numerical diffusion related to the pressure equation. We shall explain this point by studying the properties of the first-order modified equation associated to 1d Godunov like schemes which is given by

$$
\left\{\begin{array}{l}
\partial_{t} r+a_{\star} \partial_{x} u-\nu_{r} \partial_{x x}^{2} r=0  \tag{10}\\
\partial_{t} u+a_{\star} \partial_{x} r-\nu_{u} \partial_{x x}^{2} u=\omega v \\
\partial_{t} v=-\omega u
\end{array}\right.
$$

where

$$
\begin{equation*}
\nu_{r}=\frac{\kappa_{r}\left|a_{\star}\right| \Delta x}{2}, \quad \nu_{u}=\frac{\kappa_{u}\left|a_{\star}\right| \Delta x}{2} \tag{11}
\end{equation*}
$$

for some mesh size $\Delta x>0$ and viscosity parameters $\kappa_{r}>0$ and $\kappa_{u}>0$ (see [8] for more details). The classical Godunov scheme corresponds to $\kappa_{r}=\kappa_{u}=1$. In the sequel, we rewrite (10) under a vector formulation

$$
\left\{\begin{array}{l}
\partial_{t} q+L_{\kappa} q=0  \tag{12}\\
q(t=0, x)=q^{0}(x)
\end{array}\right.
$$

where $L_{\kappa}$ is the following spatial differential operator

$$
L_{\kappa}=L-B_{\kappa}, \quad L q=\left(\begin{array}{c}
a_{\star} \partial_{x} u \\
a_{\star} \partial_{x} r-\omega v \\
\omega u
\end{array}\right) \quad \text { and } \quad B_{\kappa} q=\left(\begin{array}{ccc}
\nu_{r} & 0 & 0 \\
0 & \nu_{u} & 0 \\
0 & 0 & 0
\end{array}\right)\left(\begin{array}{c}
\partial_{x x}^{2} r \\
\partial_{x x}^{2} u \\
0
\end{array}\right)
$$

### 3.1 Evolution of the energy

Lemma 2. Let $q_{\kappa}$ be the solution of System (12) on $\mathbb{T}$. Then:
i. If we define the energy by $E_{\kappa}=\left\langle q_{\kappa}, q_{\kappa}\right\rangle=\|r\|^{2}+\|u\|^{2}+\|v\|^{2}$, we obtain

$$
E_{\kappa}(t \geq 0) \leq E_{\kappa}(t=0)
$$

which means that System (12) is dissipative.
ii. If we define the average of energy by $\bar{E}_{\kappa}=\|\bar{r}\|^{2}+\|\bar{u}\|^{2}+\|\bar{v}\|^{2}$ with

$$
\bar{r}(t)=\frac{1}{|\mathbb{T}|} \int_{\mathbb{T}} r(t, x) \mathrm{d} x, \quad \bar{u}(t)=\frac{1}{|\mathbb{T}|} \int_{\mathbb{T}} u(t, x) \mathrm{d} x, \quad \bar{v}(t)=\frac{1}{|\mathbb{T}|} \int_{\mathbb{T}} v(t, x) \mathrm{d} x,
$$

we obtain

$$
\bar{E}_{\kappa}(t=0)=\bar{E}_{\kappa}(t>0) .
$$

iii. Moreover, we have

$$
\forall t>0, \bar{E}_{\kappa}(0)=\bar{E}_{\kappa}(t) \leq E_{\kappa}(t) \leq E_{\kappa}(0) .
$$

Proof. We have

$$
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t}\left\|q_{\kappa}\right\|^{2}(t)=-\left\langle L q_{\kappa}, q_{\kappa}\right\rangle+\left\langle B_{\kappa} q_{\kappa}, q_{\kappa}\right\rangle .
$$

However,

$$
\left\langle L q_{\kappa}, q_{\kappa}\right\rangle=\left\langle a_{\star} \partial_{x} u, r\right\rangle+\left\langle a_{\star} \partial_{x} r-\omega v, u\right\rangle+\langle\omega u, v\rangle=0
$$

and

$$
\left\langle B_{\kappa} q_{\kappa}, q_{\kappa}\right\rangle=\left\langle\nu_{r} \frac{\partial^{2} r}{\partial x^{2}}, r\right\rangle+\left\langle\nu_{u} \frac{\partial^{2} u}{\partial x^{2}}, u\right\rangle=-\nu_{r}\left\|\partial_{x} r\right\|^{2}-\nu_{u}\left\|\partial_{x} u\right\|^{2} .
$$

For this reason, we obtain $E_{\kappa}^{\prime}(t) \leq 0$ which means that $E_{\kappa}(t \geq 0) \leq E_{\kappa}(t=0)$.
By integrating the first order modified equation over $\mathbb{T}$ and using periodic boundary conditions, we obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \bar{r}(t)=0, \quad \frac{\mathrm{~d}}{\mathrm{~d} t} \bar{u}(t)=\omega \bar{v}(t) \quad \text { and } \quad \frac{\mathrm{d}}{\mathrm{~d} t} \bar{v}(t)=-\omega \bar{u}(t)
$$

which leads to

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \bar{r}(t)^{2}=0, \quad \frac{\mathrm{~d}}{\mathrm{~d} t} \bar{u}(t)^{2}=2 \omega \bar{v}(t) \bar{u}(t) \quad \text { and } \quad \frac{\mathrm{d}}{\mathrm{~d} t} \bar{v}(t)^{2}=-2 \omega \bar{u}(t) \bar{v}(t) .
$$

As a result, we get

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left[\bar{r}(t)^{2}+\bar{u}(t)^{2}+\bar{v}(t)^{2}\right]=0
$$

which means that $\bar{E}_{\kappa}(t=0)=\bar{E}_{\kappa}(t>0)$. It is interesting to note that

$$
\begin{aligned}
E_{\star}(t) & :=\int_{\mathbb{T}}(r-\bar{r})^{2} \mathrm{~d} x+\int_{\mathbb{T}}(u-\bar{u})^{2} \mathrm{~d} x+\int_{\mathbb{T}}(v-\bar{v})^{2} \mathrm{~d} x \\
& =\int_{\mathbb{T}}\left(r^{2}+u^{2}+v^{2}\right) \mathrm{d} x-2 \bar{r} \int_{\mathbb{T}} r \mathrm{~d} x-2 \bar{u} \int_{\mathbb{T}} u \mathrm{~d} x-2 \bar{v} \int_{\mathbb{T}} v \mathrm{~d} x+\int_{\mathbb{T}}\left(\bar{r}^{2}+\bar{u}^{2}+\bar{v}^{2}\right) \mathrm{d} x \\
& =\int_{\mathbb{T}}\left(r^{2}+u^{2}+v^{2}\right) \mathrm{d} x-\int_{\mathbb{T}}\left(\bar{r}^{2}+\bar{u}^{2}+\bar{v}^{2}\right) \mathrm{d} x=E_{\kappa}(t)-\bar{E}_{\kappa}(t) .
\end{aligned}
$$

Therefore, we obtain $E_{\star}^{\prime}(t)=E_{\kappa}^{\prime}(t) \leq 0$ and $E_{\kappa}(t) \geq \bar{E}_{\kappa}(t)\left(\right.$ since $\left.E_{\star}(t) \geq 0\right)$.

### 3.2 Structure of the kernel of the modified equation

Interestingly, the structure of the kernel ker $L_{\kappa}$ is deeply related to the value of $\kappa_{r}$. Indeed, we have:

## Lemma 3.

i. When $\kappa_{r}=0$, the subspace $\mathcal{E}_{\omega \neq 0}$ is also the kernel of the modified equation

$$
\operatorname{ker} L_{\kappa_{r}=0}=\mathcal{E}_{\omega \neq 0} .
$$

Moreover, $\mathcal{E}_{\omega \neq 0}^{\perp}$ is invariant by the modified equation.
ii. When $\kappa_{r} \neq 0$, the subspace $\mathcal{E}_{\omega \neq 0}$ is not invariant for the modified equation since

$$
\text { ker } L_{\kappa_{r} \neq 0}=\{q:=(r, u, v) \mid r=\text { const }, u=0, v=0\} \subsetneq \mathcal{E}_{\omega \neq 0} .
$$

Proof. With $\kappa_{r}=0$, it is easy to see that

$$
\operatorname{ker} L_{\kappa_{r}=0}=\mathcal{E}_{\omega \neq 0} .
$$

As for the orthogonal space, the proof of Prop. 1 (ii) stands for $\kappa_{r}=0$.
We now focus on the case $\kappa_{r} \neq 0$. Let us suppose that $q=(r, u, v) \in \operatorname{ker} L_{\kappa}$. As $u=0$, we have $a_{\star} \partial_{x} r-\omega v=0$. Then, from $L_{\kappa} q=0$, we deduce

$$
0=\left\langle L_{k} q, q\right\rangle=\nu_{r}\left\|\partial_{x} r\right\|^{2}
$$

which implies that $\partial_{x} r=0$ or equivalently $r$ is a constant. This leads to $v=0$ and $q=($ const, 0,0$)$.
The result in Lemma 3 indicates that the classical Godunov scheme ( $\kappa_{r}=1$ ) does not capture all states $q \in \mathcal{E}_{\omega \neq 0}$ because of the fact that $\operatorname{ker} L_{\kappa_{r}=1}$ is a proper subset of $\mathcal{E}_{\omega \neq 0}$. This gives rise to the loss of invariance. However, when the numerical viscosity on the pressure vanishes ( $\kappa_{r}=0$ ), we recover all states $q \in \mathcal{E}_{\omega \neq 0}$.

### 3.3 Behaviour of the solution of the modified equation

We recall that $M$ is a small parameter. Let us introduce the following definitions:
Definition 1. $A$ state $q^{0}$ is said to be well-prepared if $\left\|q^{0}-\mathbb{P} q^{0}\right\|=\mathcal{O}(M)$, where $\mathbb{P}$ is defined by (9).
Definition 2. The solution $q_{\kappa}$ of System (12) is said to be accurate at low Froude number at any time if:

$$
\forall C_{1}>0, \exists C_{2}>0,\left\|q^{0}-\mathbb{P}^{0}\right\| \leq C_{1} M \Longrightarrow \forall t \geq 0,\left\|q_{\kappa}-\mathbb{P} q^{0}\right\|(t) \leq C_{2} M,
$$

where $C_{2}$ is a positive parameter that does not depend on $M$.
Definition 3. The solution $q_{\kappa}$ of System (12) is said to be accurate at low Froude number locally in time $i f:$

$$
\forall C_{1}>0, \forall C_{2}>0: C_{2}=\mathcal{O}(1), \exists C_{3}>0,\left\|q^{0}-\mathbb{P} q^{0}\right\| \leq C_{1} M \Longrightarrow \forall t \leq C_{2},\left\|q_{\kappa}-\mathbb{P} q^{0}\right\|(t) \leq C_{3} M
$$

where $C_{3}=\mathcal{O}(1)$.
Remark 2. We notice that if the solution is accurate at low Froude number, it is free of spurious acoustic waves (refer to [8] for more details).

We have the following result:
Theorem 1. Let $q_{\kappa}$ is the solution of System (12). Then:
i. When $\kappa_{r}=0$, the solution is accurate at low Froude number at any time. Moreover, it satisfies $\left\|q_{\kappa}-\mathbb{P} q^{0}\right\|(t) \leq\left\|q^{0}-\mathbb{P} q^{0}\right\|$.
ii. When $\kappa_{r}=\mathcal{O}(M)$, the solution is accurate at low Froude number locally in time.
iii. When $\kappa_{r}=\mathcal{O}(1)$, the solution is accurate at low Froude number locally in time if $\Delta x=\mathcal{O}(M)$.

Remark 3. From Point (iii), we can state that for $\kappa_{r}=\mathcal{O}(1)$, it is enough to consider a very fine mesh to obtain accurate results. We shall see in the sequel that we actually need to consider fine meshes which is a strong restriction from the computational point of view.

Proof. Let $q_{\kappa}^{a}$ be the solution of

$$
\left\{\begin{array}{l}
\partial_{t} q+L_{\kappa} q=0 \\
q(t=0, x)=\mathbb{P} q^{0}(x)
\end{array}\right.
$$

and $q_{\kappa}^{b}$ be the solution of

$$
\left\{\begin{array}{l}
\partial_{t} q+L_{\kappa} q=0 \\
q(t=0, x)=q^{0}(x)-\mathbb{P}^{0}(x)
\end{array}\right.
$$

Then by linearity the solution of (12) is $q_{\kappa}=q_{\kappa}^{a}+q_{\kappa}^{b}$. If we suppose that $\left\|q^{0}-\mathbb{P} q^{0}\right\|=C_{1} M$, then by applying Lemma 2, we obtain

$$
\begin{equation*}
\left\|q_{\kappa}^{b}\right\|(t) \leq\left\|q_{\kappa}^{b}\right\|(0)=\left\|q^{0}-\mathbb{P}^{0}\right\|=C_{1} M \tag{13}
\end{equation*}
$$

We also notice that

$$
\begin{equation*}
\left\|q_{\kappa}-\mathbb{P} q^{0}\right\|(t)=\left\|q_{\kappa}^{a}+q_{\kappa}^{b}-\mathbb{P}^{0}\right\|(t) \leq\left\|q_{\kappa}^{a}-\mathbb{P} q^{0}\right\|(t)+\left\|q_{\kappa}^{b}\right\|(t) . \tag{14}
\end{equation*}
$$

If $\kappa_{r}=0$, then $q_{\kappa}^{a}=\mathbb{P} q^{0}$ according to Lemma 3 (i). This proves Point (i).
As for Points (ii) and (iii), we set $\hat{q}^{0}=\left(\hat{r}^{0}, \hat{u}^{0}, \hat{v}^{0}\right):=\mathbb{P} q^{0}$ and $q_{\kappa}^{a}=\left(r_{\kappa}^{a}, u_{\kappa}^{a}, v_{\kappa}^{a}\right)$. Then, we obtain

$$
\left\{\begin{array}{l}
\partial_{t}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)+a_{\star} \partial_{x}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)-\nu_{r} \partial_{x x}^{2}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)+a_{\star} \partial_{x} \hat{u}^{0}-\nu_{r} \partial_{x x}^{2} \hat{r}^{0}=0,  \tag{15}\\
\partial_{t}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)+a_{\star} \partial_{x}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)-\nu_{u} \partial_{x x}^{2}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)+a_{\star} \partial_{x} r^{0}-\nu_{u} \partial_{x x}^{2} \hat{u}^{0}=\omega\left(v_{\kappa}^{a}-\hat{v}^{0}\right)+\omega \hat{v}^{0}, \\
\partial_{t}\left(v_{\kappa}^{a}-\hat{v}^{0}\right)+\omega\left(u_{\kappa}^{a}-\hat{u}^{0}\right)+\omega \hat{u}^{0}=0 .
\end{array}\right.
$$

On the other hand, since $\mathbb{P} q^{0} \in \mathcal{E}_{\omega \neq 0}$, we have that $\hat{u}^{0}=0$ and $a_{\star} \partial_{x} \hat{r}^{0}=\omega \hat{v}^{0}$. Therefore, (15) reduces to

$$
\left\{\begin{array}{l}
\partial_{t}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)+a_{\star} \partial_{x}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)-\nu_{r} \partial_{x x}^{2}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)-\nu_{r} \partial_{x x}^{2} \hat{r}^{0}=0,  \tag{16}\\
\partial_{t}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)+a_{\star} \partial_{x}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)-\nu_{u} \partial_{x x}^{2}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)=\omega\left(v_{\kappa}^{a}-\hat{v}^{0}\right), \\
\partial_{t}\left(v_{\kappa}^{a}-\hat{v}^{0}\right)=-\omega\left(u_{\kappa}^{a}-\hat{u}^{0}\right) .
\end{array}\right.
$$

Multiplying Equation (16) by $q_{\kappa}^{a}-\hat{q}^{0}$, integrating over $\mathbb{T}$ and using periodic boundary conditions, we obtain

$$
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t}\left\|q_{\kappa}^{a}-\mathbb{P} q^{0}\right\|^{2}=-\nu_{r}\left\|\partial_{x}\left(r_{\kappa}^{a}-\hat{r}^{0}\right)\right\|^{2}-\nu_{u}\left\|\partial_{x}\left(u_{\kappa}^{a}-\hat{u}^{0}\right)\right\|^{2}+\nu_{r}\left\langle\partial_{x x}^{2} \hat{r}^{0}, r_{\kappa}^{a}-\hat{r}^{0}\right\rangle
$$

which yields

$$
\frac{1}{2} \frac{\mathrm{~d}}{\mathrm{~d} t}\left\|q_{\kappa}^{a}-\mathbb{P}^{0}\right\|^{2} \leq \nu_{r}\left\|\partial_{x x}^{2} \hat{r}^{0}\right\| \cdot\left\|r_{\kappa}^{a}-\hat{r}^{0}\right\| \leq \nu_{r}\left\|\partial_{x x}^{2} \hat{r}^{0}\right\| \cdot\left\|q_{\kappa}^{a}-\mathbb{P} q^{0}\right\|
$$

This leads to

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left\|q_{\kappa}^{a}-\mathbb{P}^{0}\right\| \leq \nu_{r}\left\|\partial_{x x}^{2} \hat{r}^{0}\right\|
$$

We deduce from the latter inequality that

$$
\begin{equation*}
\left\|q_{k}^{a}-\mathbb{P}^{0} q^{0}\right\|(t) \leq \nu_{r} t\left\|\partial_{x x}^{2} \hat{r}^{0}\right\| \tag{17}
\end{equation*}
$$

since $q_{\kappa}^{a}(0)=\mathbb{P} q^{0}$. From (13), (14) and (17), we infer

$$
\left\|q_{\kappa}-\mathbb{P}^{0}\right\|(t) \leq C_{1} M+\nu_{r} t\left\|\partial_{x x}^{2} \hat{r}^{0}\right\|
$$

With $\nu_{r}=\frac{\kappa_{r}\left|a_{\star}\right| \Delta x}{2}$, we deduce Points (ii) and (iii) respectively for $\kappa_{r}=\mathcal{O}(M)$ and $\Delta x=\mathcal{O}(M)$.

### 3.4 Fourier analysis

To go further in the study of the accuracy of the numerical scheme, we perform a Fourier analysis to investigate diffusion and dispersion effects. Let us consider functions of the form

$$
\begin{equation*}
q(t, x)=e^{i(\tau t+k x)} \hat{q} \tag{18}
\end{equation*}
$$

where $k$ is the wave number and $\tau$ is the frequency of the wave. These functions can be solutions to the modified equation only under a dispersion relation between $\tau$ and $k$ which is commonly written as $\tau=\tau(k)$. In general, this relation lies in the complex set: the real part $\Re(\tau)$ and the imaginary part $\Im(\tau)$ indicate respectively propagation and decay of Fourier modes.

Given a wave number $k$, we only consider mesh sizes satisfying

$$
\begin{equation*}
k<\frac{\pi}{\Delta x} \tag{19}
\end{equation*}
$$

so that the associated wave is captured by the scheme.
Functions (18) are solutions to the modified equation (10) if

$$
i \tau \hat{q}+A \hat{q}=0, \quad \text { where } \quad A\left(k, \nu_{r}, \nu_{u}, a_{\star}, \omega\right)=\left(\begin{array}{ccc}
\nu_{r} k^{2} & a_{\star} i k & 0  \tag{20}\\
a_{\star} i k & \nu_{u} k^{2} & -\omega \\
0 & \omega & 0
\end{array}\right) .
$$

This means that $-i \tau$ is an eigenvalue of $A$. We shall denote by $\lambda$ the eigenvalues of $A$ in the sequel. Hence the decay of Fourier mode $k$ corresponds to $\Re(\lambda) \geq 0$.

Proposition 3. Under Hypothesis (19), the damping of Fourier modes is parametrised by $\kappa_{r}$ as follows.
i. When $\kappa_{r}=0$, the wave associated to the kernel of the wave operator is preserved $(\lambda=0)$.
ii. When $\kappa_{r}=\mathcal{O}(M)$, the wave resulting from $\lambda\left(\kappa_{r}=0\right)=0$ is damped at an $\mathcal{O}(M)$ speed.
iii. When $\kappa_{r}=\mathcal{O}(1)$ and $\Delta x=\mathcal{O}(1)$, all Fourier modes are strongly damped at an $\mathcal{O}(1)$ speed.

Proof. The linear system (20) reads in terms of eigenvalues $\lambda$

$$
\begin{align*}
\nu_{r} k^{2} r+i k a_{\star} u & =\lambda r,  \tag{21a}\\
i k a_{\star} r+\nu_{u} k^{2} u-\omega v & =\lambda u,  \tag{21b}\\
\omega u & =\lambda v . \tag{21c}
\end{align*}
$$

The characteristic polynomial of Matrix $A$ is

$$
\begin{equation*}
\chi\left(\lambda, \nu_{r}\right):=\lambda^{3}-k^{2}\left(\nu_{r}+\nu_{u}\right) \lambda^{2}+\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right) \lambda-k^{2} \omega^{2} \nu_{r}=0 . \tag{22}
\end{equation*}
$$

It is a third order polynomial whose highest order coefficient is equal to one. It thus has either one real root and two complex conjugate roots (denoted respectively by $\lambda_{0}, \lambda_{c}$ and $\bar{\lambda}_{c}$ ) or three real roots (denoted respectively by $\lambda_{0}, \lambda_{+}$and $\lambda_{-}$).

It is possible to determine its three roots when $\kappa_{r}=0$ :

$$
\begin{gathered}
\lambda_{0}\left(\nu_{r}=0\right)=0, \\
\lambda_{c}\left(\nu_{r}=0\right)=\frac{1}{2}\left[k^{2} \nu_{u}+i \sqrt{4\left(\omega^{2}+k^{2} a_{\star}^{2}\right)-k^{4} \nu_{u}^{2}}\right] .
\end{gathered}
$$

We mention that the term under the square root is actually positive under Hyp. (19) (see (11) for the definition of $\nu_{u}$ ). Point (i) is proven.

We remark that $\partial_{\lambda} \chi$ does not vanish as soon as

$$
\begin{equation*}
k^{2} \Delta x^{2}\left(\left(\kappa_{r}-\kappa_{u}\right)^{2}+\kappa_{r} \kappa_{u}\right)<12\left(1+\frac{\omega^{2}}{a_{\star}^{2} k^{2}}\right) . \tag{23}
\end{equation*}
$$

Due to Hyp. (19), this inequality always holds for $\kappa_{r}$ and $\kappa_{u}$ in $[0,1]$. Hence by means of the implicit function theorem, we can define a function $\nu_{r} \mapsto \lambda_{0}\left(\nu_{r}\right)$ for $\nu_{r}$ small enough. Since coefficients multiplying $\lambda^{k}$ in (22) are affine functions in $\nu_{r}$, we infer that $\lambda_{0}$ is continuous and analytic with respect to $\nu_{r}$ [14]. This shows that

$$
\lambda_{0}\left(\nu_{r}\right) \underset{\nu_{r} \rightarrow 0}{\sim} \lambda_{0}^{\prime}\left(\nu_{r}=0\right) \nu_{r}=-\frac{\partial_{\nu_{r}} \chi(0,0)}{\partial_{\lambda} \chi(0,0)} \nu_{r}=\frac{k^{2} \omega^{2}}{k^{2} a_{\star}^{2}+\omega^{2}} \nu_{r} .
$$

In particular, we deduce that if $\kappa_{r}=\mathcal{O}(M)$, then $\lambda_{0}\left(\nu_{r}\right)=\mathcal{O}(M)$. This proves Point (ii).
Let us now provide other properties of the eigenvalues. We substitute (21c) into (21b) and then multiply (21a) by $\bar{r}$ and (21b) by $\bar{u}$ to obtain

$$
\begin{equation*}
\frac{1}{\lambda} \omega^{2}|u|^{2}+\lambda\left(|r|^{2}+|u|^{2}\right)=k^{2}\left(\nu_{r}|r|^{2}+\nu_{u}|u|^{2}\right)+i k a_{\star}(u \bar{r}+r \bar{u}) . \tag{24}
\end{equation*}
$$

On the one hand, the real part of (24)

$$
\Re(\lambda)\left[\frac{\omega^{2}|u|^{2}}{|\lambda|^{2}}+|r|^{2}+|u|^{2}\right]=k^{2}\left(\nu_{r}|r|^{2}+\nu_{u}|u|^{2}\right),
$$

shows that all eigenvalues have positive real parts (unless $k=0$ for which eigenvalues are pure imaginary), which ensures the decay for all Fourier modes.

The three roots of (22) satisfy

$$
\begin{align*}
\lambda_{1}+\lambda_{2}+\lambda_{3} & =k^{2}\left(\nu_{r}+\nu_{u}\right),  \tag{25a}\\
\lambda_{1} \lambda_{2}+\left(\lambda_{1}+\lambda_{2}\right) \lambda_{3} & =\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u},  \tag{25b}\\
\lambda_{1} \lambda_{2} \lambda_{3} & =k^{2} \omega^{2} \nu_{r} . \tag{25c}
\end{align*}
$$

Substituting $\lambda_{1} \lambda_{2}$ from (25c) into (25b), we get

$$
\begin{equation*}
\left(\lambda_{1}+\lambda_{2}\right) \lambda_{3}+\frac{k^{2} \omega^{2} \nu_{r}}{\lambda_{3}}=\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u} . \tag{26}
\end{equation*}
$$

Let us first focus on the case of a single real eigenvalue: we take $\lambda_{1}=\lambda_{c}, \lambda_{2}=\bar{\lambda}_{c}$ and $\lambda_{3}=\lambda_{0}$. Eq. (26) yields

$$
\begin{equation*}
\lambda_{0} \geq \frac{k^{2} \omega^{2} \nu_{r}}{\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}} \tag{27}
\end{equation*}
$$

since it has been proven that $\Re\left(\lambda_{c}\right) \geq 0$.
We also notice that $\chi\left(0, \nu_{r}\right)=-k^{2} \omega^{2} \nu_{r}<0$ and $\chi\left(k^{2} \nu_{r}, \nu_{r}\right)=k^{4} a_{\star}^{2} \nu_{r}>0$. Hence, since there is a single real eigenvalue, this implies that $\lambda_{0} \leq k^{2} \nu_{r}$ and we have

$$
\begin{equation*}
\frac{\omega^{2}}{\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}} k^{2} \nu_{r} \leq \lambda_{0} \leq k^{2} \nu_{r} . \tag{28}
\end{equation*}
$$

As for the complex conjugate roots, we get from (25a) and (26) that $\mu:=2 \Re\left(\lambda_{c}\right)$ verifies

$$
\begin{equation*}
f(\mu):=\mu^{2}-k^{2}\left(\nu_{r}+\nu_{u}\right) \mu-\frac{k^{2} \omega^{2} \nu_{r}}{k^{2}\left(\nu_{r}+\nu_{u}\right)-\mu}+\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}=0 . \tag{29}
\end{equation*}
$$

We remark that $f(\mu) \geq g(\mu)$ where

$$
\begin{equation*}
g(\mu):=-k^{2}\left(\nu_{r}+\nu_{u}\right) \mu-\frac{k^{2} \omega^{2} \nu_{r}}{k^{2}\left(\nu_{r}+\nu_{u}\right)-\mu}+\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u} . \tag{30}
\end{equation*}
$$

Since $f(0)=g(0)>0$, this implies that any root $\mu$ of (29) is larger than the smallest positive root of (30). Equation $g(\mu)=0$ can be written as

$$
\begin{equation*}
k^{2}\left(\nu_{r}+\nu_{u}\right) \mu^{2}-\left[k^{4}\left(\nu_{r}+\nu_{u}\right)^{2}+\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right)\right] \mu+\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right) k^{2}\left(\nu_{r}+\nu_{u}\right)-k^{2} \omega^{2} \nu_{r}=0 \tag{31}
\end{equation*}
$$

Due to the fact that

$$
\Delta=\left[k^{4}\left(\nu_{r}+\nu_{u}\right)^{2}-\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right)\right]^{2}+4 k^{4}\left(\nu_{r}+\nu_{u}\right) \nu_{r} \omega^{2}>0
$$

Equation (31) has two real positive solutions so that

$$
\begin{equation*}
2 \Re\left(\lambda_{c}\right) \geq \frac{\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right) k^{2}\left(\nu_{r}+\nu_{u}\right)-k^{2} \omega^{2} \nu_{r}}{k^{4}\left(\nu_{r}+\nu_{u}\right)^{2}+\left(\omega^{2}+k^{2} a_{\star}^{2}+k^{4} \nu_{r} \nu_{u}\right)} \tag{32}
\end{equation*}
$$

In the case of three real roots, (27) holds for each of them by symmetry as they are all positive. Lower bounds (27) and (32) ensure that real parts of all eigenvalues are of order 1 when $\nu_{r}$ is of order 1. This proves Point (iii).

## 4 Analysis of fully discrete Godunov schemes

There are two main possible time strategies for Godunov type schemes applied to the linear wave equation with Coriolis source term. The first one is a classical splitting discretisation where one deals with the problem without source term in a first step and then the Coriolis source term is considered in a second step, which then consists in solving an ordinary differential equation. It is well known that this splitting strategy is not well adapted to preserve stationary states and then to compute small perturbations around them [3, 12], see also Appendix A. Thus we focus on the analysis of the second strategy that consists in computing acoustic and Coriolis effects in a single step. As a matter of fact, there are many ways to take into account the Coriolis source term. For example, we can discretise this term using explicit, implicit and even Crank-Nicolson strategies. Hence, we introduce two new parameters $\theta_{1}$ and $\theta_{2}$ to parametrise the strategy.

### 4.1 Study of the discrete kernel of the one step Godunov scheme

We consider a homogeneous cartesian mesh $\left(x_{i}\right)_{1 \leq i \leq N}$. The one step fully discrete Godunov scheme is given by

$$
\left\{\begin{array}{l}
\frac{r_{i}^{n+1}-r_{i}^{n}}{\Delta t}+a_{\star} \frac{u_{i+1}^{n}-u_{i-1}^{n}}{2 \Delta x}-\nu_{r} \frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}=0  \tag{33}\\
\frac{u_{i}^{n+1}-u_{i}^{n}}{\Delta t}+a_{\star} \frac{r_{i+1}^{n}-r_{i-1}^{n}}{2 \Delta x}-\nu_{u} \frac{u_{i+1}^{n}-2 u_{i}^{n}+u_{i-1}^{n}}{\Delta x^{2}}=\omega\left[\theta_{1} v_{i}^{n}+\left(1-\theta_{1}\right) v_{i}^{n+1}\right] \\
\frac{v_{i}^{n+1}-v_{i}^{n}}{\Delta t}=-\omega\left[\theta_{2} u_{i}^{n}+\left(1-\theta_{2}\right) u_{i}^{n+1}\right]
\end{array}\right.
$$

for $i \in\{1, \ldots, N\}$ and $0 \leq \theta_{1}, \theta_{2} \leq 1$. Periodic boundary conditions read

$$
\begin{equation*}
q_{0}^{n+1}=q_{N}^{n+1}, q_{N+1}^{n+1}=q_{1}^{n+1} \tag{34}
\end{equation*}
$$

We now investigate the kernel of the fully discrete one step scheme. It is strongly related to the value of the numerical viscosity $\kappa_{r}$. In particular we have the following result:

## Lemma 4.

i. When $\kappa_{r}=0$, the kernel of the one step scheme is

$$
\mathcal{E}_{\omega \neq 0}^{h}:=\operatorname{ker} L_{\kappa_{r}=0, h}=\left\{q=(r, u, v) \in \mathbb{R}^{3 N} \mid u_{i}=0, \frac{a_{\star}}{2 \Delta x}\left(r_{i+1}-r_{i-1}\right)=\omega v_{i}\right\}
$$

ii. When $\kappa_{r} \neq 0$, the kernel of the one step scheme is

$$
\operatorname{ker} L_{\kappa_{r} \neq 0, h}=\left\{q=(r, u, v) \in \mathbb{R}^{3 N} \mid \exists C \in \mathbb{R}: r_{i}=C, u_{i}=0, v_{i}=0\right\}
$$

Proof. A stationary state verifies $r_{i}^{n+1}=r_{i}^{n}, u_{i}^{n+1}=u_{i}^{n}$ and $v_{i}^{n+1}=v_{i}^{n}$. Therefore, we easily obtain from (33) that

$$
\left\{\begin{array}{l}
a_{\star} \frac{u_{i+1}^{n}-u_{i-1}^{n}}{2 \Delta x}-\nu_{r} \frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}=0  \tag{35a}\\
a_{\star} \frac{r_{i+1}^{n}-r_{i-1}^{n}}{2 \Delta x}-\nu_{u} \frac{u_{i+1}^{n}-2 u_{i}^{n}+u_{i-1}^{n}}{\Delta x^{2}}=\omega v_{i}^{n} \\
0=-\omega u_{i}^{n}
\end{array}\right.
$$

Point (i) is straightforward: we get from (35c) that $u_{i}^{n}=0$, and then (35a) is trivially satisfied since $\kappa_{r}=0$. Then, (35b) yields that

$$
\frac{a_{\star}}{2 \Delta x}\left(r_{i+1}^{n}-r_{i-1}^{n}\right)=\omega v_{i}^{n}
$$

Now we consider the case $\kappa_{r} \neq 0$. According to (35c), $u_{i}^{n}=0$ for all $i$. Together with (35a) and $\nu_{r} \neq 0$, we get $r_{i+1}^{n}-r_{i}^{n}=r_{i}^{n}-r_{i-1}^{n}$. By induction we get $r_{N+1}^{n}-r_{N}^{n}=r_{N}^{n}-r_{N-1}^{n}=\ldots=r_{2}^{n}-r_{1}^{n}=r_{1}^{n}-r_{0}^{n}=c$ where $c$ is a constant. This implies $r_{N}^{n}=r_{0}^{n}+N c$. On the other hand, periodic conditions require to have $r_{N}^{n}=r_{0}^{n}$. Therefore, we get $c=0$ and $r_{i}^{n}=$ constant. This leads to $v_{i}^{n}=0$ by using (35b). Point (ii) is proven.

### 4.2 Stability of the discrete one step Godunov scheme

For $0 \leq \theta_{1}, \theta_{2} \leq 1$, let us denote

$$
\Theta_{1}=1-\theta_{1}-\theta_{2}, \quad \Theta_{2}=\theta_{1} \theta_{2}+\left(1-\theta_{1}\right)\left(1-\theta_{2}\right) \in[0,1], \quad \Theta_{3}=\left(1-2 \theta_{1}\right)\left(1-2 \theta_{2}\right) \in[-1,1]
$$

Lemma 5. For $\kappa_{r}=0$ and $\kappa_{u}>0$, we have:
i. When $\theta_{1}+\theta_{2}>1$, the one step scheme (33) is unstable.
ii. When $\theta_{1}+\theta_{2} \leq 1$, we consider two cases:
(a) If $\frac{\kappa_{u}^{2} a_{\star}^{2}}{\omega^{2} \Delta x^{2}} \leq \Theta_{3}$, the one step scheme (33) is stable provided that

$$
\begin{equation*}
\Delta t \leq \Delta t_{a}:=\frac{\kappa_{u} \Delta x}{2\left|a_{\star}\right|} \frac{1}{\left(1-\frac{\omega \Delta x}{\left|a_{\star}\right|} \sqrt{\Theta_{1}}\right)_{+}} \tag{36a}
\end{equation*}
$$

(b) If $\frac{\kappa_{u}^{2} a_{\star}^{2}}{\omega^{2} \Delta x^{2}}>\Theta_{3}$, the one step scheme (33) is stable provided that

$$
\Delta t \leq \min \left\{\Delta t_{a}, \Delta t_{b}\right\} \quad \text { where } \quad \Delta t_{b}:=\frac{\Delta x}{\kappa_{u}\left|a_{\star}\right|} \times \begin{cases}\frac{2 \kappa_{u}^{2} a_{\star}^{2}}{\omega^{2} \Delta x^{2} \Theta_{3}}\left[1-\sqrt{1-\frac{\omega^{2} \Delta x^{2}}{\kappa_{u}^{2} a_{\star}^{2}} \Theta_{3}}\right], & \text { if } \Theta_{3} \neq 0  \tag{36b}\\ 1, & \text { otherwise } .\end{cases}
$$



Figure 1: Region of stability condition.

Remark 4. The standard CFL condition for the homogeneous case ( $\omega=0$ ) reads [8]

$$
\Delta t \leq \Delta t_{0}:=\frac{\Delta x}{\left|a_{\star}\right|} \min \left\{\frac{\kappa_{u}}{2}, \frac{1}{\kappa_{u}}\right\} .
$$

Inequality (36a) clearly shows that taking Coriolis forces into account requires a less restrictive CFL condition. It is also the case for (36b) when $\Theta_{3} \geq 0$ thanks to the convexity of the function $x \mapsto 1-\sqrt{1-x}$. We also notice that for the Crank-Nicolson scheme $\theta_{1}=\theta_{2}=\frac{1}{2}$, we recover the standard bound $\Delta t_{0}$.
Remark 5. An asymptotic expansion for $\Delta x \ll 1$ in the bound $\Delta t_{a}$ and $\Delta t_{b}$ in (36a-36b) yields

$$
\Delta t_{a}=\frac{\kappa_{u} \Delta x}{2\left|a_{\star}\right|}+\mathcal{O}\left(\Delta x^{2}\right), \quad \Delta t_{b}=\frac{\Delta x}{\kappa_{u}\left|a_{\star}\right|}+\mathcal{O}\left(\Delta x^{3}\right)
$$

and then one still recovers the classical bound $\Delta t_{0}$ for the homogeneous problem.
Remark 6. For large values of the Coriolis parameter $\omega$, the constraint (36a) is always satisfied ( $\Delta t_{a}=+\infty$ ) while for the second constraint (36b), it depends on the sign of $\Theta_{3}$ :

- If $\Theta_{3} \geq 0$, there is no constraint upon $\Delta t$ for $\omega$ large enough;
- If $\Theta_{3}<0$, the asymptotic bound reads

$$
\Delta t_{b} \approx \frac{2}{\omega} .
$$

We then recover the standard stability condition for the ODE system solved by means of a $\theta$-scheme (44b). Remark 7. Figure 1 specifies the stability area. In the red zone, the scheme is unstable according to Point (i). In the green zone, the scheme is stable under a CFL-like constraint (characterised by $\Delta t_{a}$ or $\left.\min \left(\Delta t_{a}, \Delta t_{b}\right)\right)$ that is less restrictive than the homogeneous bound $\Delta t_{0}$ while in the blue zone, the scheme is stable provided $\Delta t$ is smaller than $\min \left(\Delta t_{a}, \Delta t_{b}\right) \leq \Delta t_{0}$.

Proof. We perform a Von Neumann analysis to investigate the stability condition for Scheme (33). Let us denote

$$
\sigma=\frac{\Delta t}{\Delta x}, \gamma=\omega \Delta t \text { and } s=\sin \left(\frac{k \Delta x}{2}\right) .
$$

We now substitute

$$
q_{j}^{n}=\left(\begin{array}{l}
r_{j}^{n} \\
u_{j}^{n} \\
v_{j}^{n}
\end{array}\right)=\left(\begin{array}{l}
R_{n} \\
U_{n} \\
V_{n}
\end{array}\right) e^{i k j \Delta x}
$$

into (33) in order to obtain

$$
\begin{equation*}
A q_{j}^{n+1}=B q_{j}^{n} \tag{37}
\end{equation*}
$$

where the matrices $A$ and $B$ are given by

$$
A=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & -\left(1-\theta_{1}\right) \gamma \\
0 & \left(1-\theta_{2}\right) \gamma & 1
\end{array}\right) \quad \text { and } \quad B=\left(\begin{array}{ccc}
1-2 \kappa_{r}\left|a_{\star}\right| \sigma s^{2} & -a_{\star} \sigma i \sin (k \Delta x) & 0 \\
-a_{\star} \sigma i \sin (k \Delta x) & 1-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2} & \theta_{1} \gamma \\
0 & -\theta_{2} \gamma & 1
\end{array}\right)
$$

In addition, we have

$$
A^{-1}=\frac{1}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\left(\begin{array}{ccc}
\Lambda\left(\theta_{1}, \theta_{2}\right) & 0 & 0 \\
0 & 1 & \gamma\left(1-\theta_{1}\right) \\
0 & -\gamma\left(1-\theta_{2}\right) & 1
\end{array}\right)
$$

with

$$
\begin{equation*}
\Lambda\left(\theta_{1}, \theta_{2}\right)=1+\gamma^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right) \tag{38}
\end{equation*}
$$

Therefore, we can rewrite (37) as the following equation

$$
q_{j}^{n+1}=C q_{j}^{n}
$$

where the amplification matrix $C$ is given by

$$
C=A^{-1} B=\frac{1}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\left(\begin{array}{ccc}
\left(1-2 \kappa_{r}\left|a_{\star}\right| \sigma s^{2}\right) \Lambda\left(\theta_{1}, \theta_{2}\right) & -a_{\star} \sigma i \sin (k \Delta x) \Lambda\left(\theta_{1}, \theta_{2}\right) & 0  \tag{39}\\
-a_{\star} \sigma i \sin (k \Delta x) & 1-\gamma^{2} \theta_{2}\left(1-\theta_{1}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2} & \gamma \\
\gamma\left(1-\theta_{2}\right) a_{\star} \sigma i \sin (k \Delta x) & -\gamma\left[1-\left(1-\theta_{2}\right) 2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}\right] & 1-\gamma^{2} \theta_{1}\left(1-\theta_{2}\right)
\end{array}\right)
$$

whose characteristic polynomial will be denoted by $\mathcal{P}(\lambda)$. We now consider the modes which are constant in space $(k=0)$. In this case, the amplification matrix in $(u, v)$ is given by

$$
\frac{1}{1+\gamma^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right)}\left(\begin{array}{cc}
1-\gamma^{2} \theta_{2}\left(1-\theta_{1}\right) & \gamma \\
-\gamma & 1-\gamma^{2} \theta_{1}\left(1-\theta_{2}\right)
\end{array}\right)
$$

Therefore the characteristic equation $\mathcal{P}(\lambda)=0$ reduces to

$$
\begin{equation*}
\lambda^{2}-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \lambda+\frac{1+\gamma^{2} \theta_{1} \theta_{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}=0 \tag{40}
\end{equation*}
$$

and the condition $\left|\lambda_{1} \lambda_{2}\right| \leq 1$ is equivalent to

$$
1+\gamma^{2} \theta_{1} \theta_{2} \leq 1+\gamma^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right)
$$

that is fulfilled if and only if

$$
\gamma^{2}\left[\left(\theta_{1}+\theta_{2}\right)-1\right] \leq 0
$$

which leads to the condition $\theta_{1}+\theta_{2} \leq 1$. This proves Point (i).
Now we consider the case of interest $\kappa_{r}=0$ (c.f. Lemma 4). The characteristic polynomial $\mathcal{P}(\lambda)$ reduces to
$\mathcal{P}_{0}(\lambda)=(1-\lambda)\left[\lambda^{2}-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \lambda+\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\right]$.

One root of this polynomial is $\lambda_{0}=1$ and the two others roots $\lambda_{ \pm}$are the solutions of the following second degree equation

$$
\begin{equation*}
\lambda^{2}+\xi \lambda+\zeta=0 \tag{42}
\end{equation*}
$$

with

$$
\xi=-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \quad \text { and } \quad \zeta=\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)} .
$$

In order to ensure that the roots of (42) are in the unit circle $\left(\left|\lambda_{ \pm}\right| \leq 1\right)$, the coefficients $\xi$ and $\zeta$ must satisfy

$$
|\zeta| \leq 1 \quad \text { and } \quad|\xi| \leq 1+\zeta .
$$

- Firstly, the condition $\zeta \leq 1$ is equivalent to

$$
\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \leq 1
$$

which leads to

$$
f_{1}\left(s^{2}\right):=-\gamma^{2} \Theta_{1}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right) \leq 0 .
$$

With $s$ varying in $[-1,1]$, the previous condition holds provided $\max _{[0,1]} f_{1} \leq 0$. As Function $f_{1}$ is maximal over $\mathbb{R}$ at $X_{1}:=\frac{1}{2}\left(1-\frac{\kappa_{u}}{2\left|a_{*}\right| \sigma}\right)$, we deduce that

$$
\max _{[0,1]} f_{1}= \begin{cases}f_{1}(0), & \text { if } X_{1} \leq 0 \\ f_{1}\left(X_{1}\right), & \text { otherwise }\end{cases}
$$

If $X_{1} \leq 0$ which is equivalent to $\sigma \leq \frac{\kappa_{u}}{2\left|a_{\star}\right|}$, the condition $f_{1}(0) \leq 0$ is always satisfied. If $X_{1}>0$, $f_{1}\left(X_{1}\right) \leq 0$ reads

$$
\left(\left|a_{\star}\right| \sigma-\frac{\kappa_{u}}{2}\right)^{2} \leq \gamma^{2} \Theta_{1} \Longleftrightarrow\left(\frac{\left|a_{\star}\right|}{\Delta x}-\omega \sqrt{\Theta_{1}}\right) \Delta t \leq \frac{\kappa_{u}}{2}
$$

Hence $\Delta t \leq \Delta t_{a}$.

- Next, the condition $\zeta \geq-1$ can be written as

$$
f_{2}\left(s^{2}\right):=\gamma^{2} \Theta_{2}+2\left(1-\kappa_{u}\left|a_{\star}\right| \sigma s^{2}\right)+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right) \geq 0 .
$$

We shall see below that this constraint is weaker that another one $\left(f_{3}\left(s^{2}\right) \geq 0\right)$ and needs not be taken into account.

- Let us now turn to the condition upon $\xi$. The first case $-\xi \leq 1+\zeta$ reads
$2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2} \leq 2+\gamma^{2}\left[1-\left(\theta_{1}+\theta_{2}\right)+2 \theta_{1} \theta_{2}\right]-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)$
which comes down to

$$
-\gamma^{2}-4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right) \leq 0
$$

The latter inequality always holds and does not imply an additional constraint upon $\Delta t$.

- Finally, we consider the case $\xi \leq 1+\zeta$. This leads to
$-2+\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)+2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2} \leq 2+\gamma^{2}\left[1-\left(\theta_{1}+\theta_{2}\right)+2 \theta_{1} \theta_{2}\right]-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)$.
It follows that

$$
f_{3}\left(s^{2}\right):=\gamma^{2} \Theta_{3}+4\left(1-\kappa_{u}\left|a_{\star}\right| \sigma s^{2}\right)+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right) \geq 0
$$

From $\Theta_{3}=2 \Theta_{2}-1$, we infer that $2 f_{2}\left(s^{2}\right)-f_{3}\left(s^{2}\right) \geq 0$ over [ 0,1$]$. This implies that the condition $f_{2}\left(s^{2}\right) \geq 0$ is a consequence of $f_{3}\left(s^{2}\right) \geq 0$.

Function $f_{3}$ is maximal over $\mathbb{R}$ at $X_{3}:=\frac{1}{2}\left(1-\frac{\kappa_{u}}{\left|a_{\star}\right| \sigma}\right) \leq \frac{1}{2}$. The minimum over $[0,1]$ is reached for $s^{2}=1$ and the condition $f_{3}\left(s^{2}\right) \geq 0$ reduces to

$$
0 \leq f_{3}(1)=\omega^{2} \Theta_{3} \Delta t^{2}-\frac{4 \kappa_{u}\left|a_{\star}\right|}{\Delta x} \Delta t+4=: Q_{3}(\Delta t)
$$

The resolution of the second order equation $Q_{3}(\Delta t)=0$ leads to the stability condition (36b) depending on the sign of $\omega^{2} \Theta_{3} \Delta x^{2}-\kappa_{u}^{2} a_{\star}^{2}$.

Lemma 6 (Stability of the All Froude Godunov scheme). The CFL condition (36a-36b) obtained for $\kappa_{r}=0$ still ensures the stability of the All Froude Godunov scheme, i.e. for the choice $\kappa_{r}=\mathcal{O}(M)$.

The proof is obtained by using a classical continuity argument. The key point is to prove the modulus of the eigenvalue $\lambda_{0}$ is increasing when $\kappa_{r} \rightarrow 0^{+}$. The detailed proof of Lemma 6 is given in Appendix B.

## 5 Numerical results

Let us fix the parameters $a_{\star}=1, \omega=1, M=10^{-3}$ and consider the initial condition

$$
q_{i}^{0}=\hat{q}_{i}^{0}+M \frac{\tilde{q}_{i}^{0}}{\left\|\tilde{q}_{i}^{0}\right\|} \quad \text { with } \quad \hat{q}_{i}^{0}=\left(\begin{array}{c}
\sin \left(\omega x_{i}\right) \\
0 \\
a_{\star} \cos \left(\omega x_{i}\right) \frac{\sin (\omega \Delta x)}{\omega \Delta x}
\end{array}\right) \in \mathcal{E}_{\omega \neq 0}^{h}, \tilde{q}_{i}^{0}=\left(\begin{array}{c}
a_{\star} \cos \left(\omega x_{i}\right) \frac{\sin (\omega \Delta x)}{\omega \Delta x} \\
1 \\
\sin \left(\omega x_{i}\right)
\end{array}\right) \in \mathcal{E}_{\omega \neq 0}^{h, \perp},
$$

that is close to the kernel $\mathcal{E}_{\omega \neq 0}^{h}$ (see Lemma 4) up to a perturbation of order $M$.
We solve the 1D linear wave equation (5) by means of the schemes we analyzed in the previous sections, namely the low Froude scheme (33) for $\kappa_{r}=0$, the all Froude scheme (33) for $\kappa_{r}=\mathcal{O}(M)$, and the classical Godunov scheme (33) for $\kappa_{r}=1$. In a first step, we take $\theta_{1}=1$ and $\theta_{2}=0$.


Figure 2: Evolution of $\max _{t}\left\|q_{h}-\mathbb{P} q_{h}^{0}\right\|(t)$ for $t=\mathcal{O}(1)$ when the Froude number goes to 0 for the Low Froude Godunov, the All Froude Godunov and the Classical Godunov schemes.

We observe on Figure 2 that the two schemes designed for the low Froude regime have the correct behaviour as the Froude number goes to 0 , unlike the classical Godunov scheme which is not accurate as stated before.

We now investigate the accuracy with time at a fixed Froude number. As it was stated in Theorem 1(i), we see on Figures 3(a)-(c) that the two aforementioned schemes are accurate for times $t=\mathcal{O}(1)$ since


Figure 3: Comparisons of schemes: proximity to the discrete kernel as time increases.
the numerical solutions remain close to the projection of the initial data onto the kernel (the norm of the difference is of order $10^{-3}$ ). However for large times the all Froude scheme turns out be inaccurate as the corresponding solution is moving away from the kernel. It illustrates the result from Theorem 1(ii).

Next, we now focus on Theorem 1 (iii) by means of Figures 4(a)-(b), where we see that the total deviation of the Classical Godunov scheme is of order $M$ when the mesh is sufficiently refined ( $\Delta x=\mathcal{O}(M)$ ). Note that even in this case, the behaviour of the low/all Froude Godunov schemes is better than that of the classical Godunov scheme.

In Figures 5 and 6 , we change the value $\theta_{1}$ and $\theta_{2}$ of the Low/All Froude schemes. These figures indicate that the total deviation depends on the value of $\theta_{1}+\theta_{2}$.

In the second test case, we consider the initial condition given by

$$
\left\{\begin{array}{l}
r_{i}^{0}=\chi_{\left[-\frac{1}{2}, \frac{1}{2}\right]}\left(x_{i}\right)  \tag{43}\\
u_{i}^{0}=1, \\
v_{i}^{0}=1 .
\end{array}\right.
$$

In this test, we choose $\omega=1, \Delta x=0.01$ and $a_{\star}$ such that the Rossby deformation is equal to $R_{d}:=\frac{a_{\star}}{\omega}=\Delta x$ and $\kappa_{u}=1$. In Figure 7(a), we choose $\theta_{1}=0.5$ and $\theta_{2}=0$. Hence, in this case we have $\Theta_{1}=0.5$ and $\Theta_{3}=0$ which leads to $\Delta t_{a}=\frac{0.5}{1-\sqrt{0.5}}, \Delta t_{b}=1$ and $\Delta t_{0}=0.5$. Therefore, the new time step $\Delta t=\min \left\{\Delta t_{a}, \Delta t_{b}\right\}=1$ is less restrictive than the classical time step $\Delta t_{0}=0.5$. Figure 7(a) shows that the new time step is optimal since when $\Delta t=0.999<1$ the Low Froude scheme is stable while when $\Delta t=1.001$ the Low Froude scheme is unstable.

On the other hand, if we take $\theta_{1}=0$ and $\theta_{2}=0$, then $\Theta_{1}=1$ and $\Theta_{3}=1$. Due to the fact that $\frac{\kappa_{2}^{2} a_{\star}^{2}}{\omega^{2} \Delta x^{2}} \leq \Theta_{3}$, the constraint over the time step for the Low Froude scheme is prescribed by $\Delta t_{a}$. However


Figure 4: Comparisons of schemes: proximity to the discrete kernel as time increases.


Figure 5: Comparisons of Low Froude schemes: proximity to the discrete kernel as time increases.


Figure 6: Comparisons of All Froude schemes: proximity to the discrete kernel as time increases.


Figure 7: Influence of the time step upon the Low Froude scheme
as $a_{\star}=\omega \Delta x \sqrt{\Theta_{1}}, \Delta t_{a}=+\infty$ and the Low Froude scheme is always stable without regard to the time step $\Delta t$. Figure 7(b) confirms this statement by showing that the Low Froude scheme is stable even for $\Delta t=10$.

In Figure 8, we take $\Delta t=\Delta t_{0}=0.5$ and change the value of $\theta_{1}$ and $\theta_{2}$. This figure indicates that the behaviour of the energy of the Low Froude scheme depends on the value of $\theta_{1}$ and $\theta_{2}$. The choice $\theta_{1}=\theta_{2}=\frac{1}{2}$ (Crank-Nicolson approximation for the Coriolis term) is able to preserve the energy exaclty although this choice requires a more restrictive constraint upon the time step than for $\theta_{1}, \theta_{2} \leq \frac{1}{2}$.

## 6 Conclusion

It is well known that the classical Godunov scheme applied to the linear wave equation is not accurate at low Froude number on cartesian meshes in dimension 2 [8]. In this work, we have shown that, when a Coriolis source term is involved, the classical Godunov scheme is not accurate at low Froude number even in dimension 1.

This is because the stationary space of the classical Godunov discrete operator is not a good approximation of the invariant subspace $\mathcal{E}_{\omega \neq 0}$. The loss of invariance of $\mathcal{E}_{\omega \neq 0}$ is explained by studying the associated modified equation. It is strongly related to the numerical diffusion $\kappa_{r}$ on the pressure equation. In particular when we set $\kappa_{r}=0$, the inaccuracy problem does not occur. As a result, we derived two modified schemes by decreasing the value of the numerical diffusion $\kappa_{r}$ on the pressure equation. From this, we deduce that:

- The Low Froude Godunov scheme ( $\kappa_{r}=0$ ) is accurate at low Froude number.
- The All Froude Godunov scheme $\left(\kappa_{r}=M\right)$ is accurate at low Froude number locally in time.

We then proved that both schemes are stable under suitable constraints upon the time step. These stability conditions turn out to be less restrictive than classical ones for a suitable treatment of the Coriolis source term.

In forthcoming works, we shall extend our analysis to the two-dimensional case and derive and analyse numerical schemes that are accurate and stable in this framework.


Figure 8: Comparisons of Low Froude schemes depending on the time step

## A Analysis of splitting scheme

Let us define a two-step Godunov scheme using a splitting strategy to take into account the Coriolis source term. The first step is related to the acoustic term

$$
\left\{\begin{array}{l}
r_{i}^{\star}-r_{i}^{n}+a_{\star} \Delta t \frac{u_{i+1}^{n}-u_{i-1}^{n}}{2 \Delta x}-\nu_{r} \Delta t \frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}=0  \tag{44a}\\
u_{i}^{\star}-u_{i}^{n}+a_{\star} \Delta t \frac{t_{i+1}^{n}-r_{i-1}^{n}}{2 \Delta x}-\nu_{u} \Delta t \frac{u_{i+1}^{n}-2 u_{i}^{n}+u_{i-1}^{n}}{\Delta x^{2}}=0 \\
v_{i}^{\star}-v_{i}^{n}=0
\end{array}\right.
$$

and we use a $\theta$-scheme to deal with the Coriolis term in the second step

$$
\left\{\begin{array}{l}
r_{i}^{n+1}=r_{i}^{\star},  \tag{44b}\\
u_{i}^{n+1}-u_{i}^{\star}=\omega \Delta t\left[\theta_{1} v_{i}^{\star}+\left(1-\theta_{1}\right) v_{i}^{n+1}\right] \\
v_{i}^{n+1}-v_{i}^{\star}=-\omega \Delta t\left[\theta_{2} u_{i}^{\star}+\left(1-\theta_{2}\right) u_{i}^{n+1}\right]
\end{array}\right.
$$

for $0 \leq \theta_{1}, \theta_{2} \leq 1$.

## Lemma 7.

i. For $\kappa_{r}=0$, the splitting scheme preserves steady states only if $\theta_{2}=0$.
ii. For $\kappa_{r} \neq 0$, steady states are not preserved without regard to the value of $\theta_{1}$ and $\theta_{2}$.

Proof. Let us assume that the numerical solution at time $t^{n}=n \Delta t$ belongs to the discrete kernel $\mathcal{E}_{\omega \neq 0}^{h}$

$$
\forall i \in \mathbb{Z}, \quad u_{i}^{n}=0 \quad \text { and } \quad a_{\star} \frac{r_{i+1}^{n}-r_{i-1}^{n}}{2 \Delta x}=\omega v_{i}^{n} .
$$

We shall show that at the next time step the numerical solution does not lie in the discrete kernel anymore. After the first step, we easily obtain

$$
\left\{\begin{array}{l}
r_{i}^{\star}=r_{i}^{n}+\nu_{r} \Delta t \frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}, \\
u_{i}^{\star}=u_{i}^{n}-a_{\star} \Delta t \frac{r_{i+1}^{n}-r_{i-1}^{n}}{2 \Delta x}=-\omega \Delta t v_{i}^{n}, \\
v_{i}^{\star}=v_{i}^{n} .
\end{array}\right.
$$

Then the second step leads to

$$
\left\{\begin{array}{l}
r_{i}^{n+1}=r_{i}^{n}+\nu_{r} \Delta t^{\frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}} \\
u_{i}^{n+1}=\omega \Delta t\left(1-\theta_{1}\right)\left(v_{i}^{n+1}-v_{i}^{n}\right) \\
v_{i}^{n+1}+\omega \Delta t\left(1-\theta_{2}\right) u_{i}^{n+1}=\left[1+(\omega \Delta t)^{2} \theta_{2}\right] v_{i}^{n}
\end{array}\right.
$$

As a result, we have

$$
v_{i}^{n+1}+(\omega \Delta t)^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right)\left(v_{i}^{n+1}-v_{i}^{n}\right)=\left[1+(\omega \Delta t)^{2} \theta_{2}\right] v_{i}^{n}
$$

from which it follows that

$$
\left[1+(\omega \Delta t)^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right)\right] v_{i}^{n+1}=\left[1+(\omega \Delta t)^{2}\left(1-\theta_{1}\right)\left(1-\theta_{2}\right)+(\omega \Delta t)^{2} \theta_{2}\right] v_{i}^{n}
$$

Therefore, $v_{i}^{n+1}=v_{i}^{n}\left(\right.$ and $\left.u_{i}^{n+1}=0\right)$ iff $\theta_{2}=0$. The kernel is recovered if $\kappa_{r}=0$ as $r_{i}^{n+1}=r_{i}^{n}$.
Let us now note that the choice $\theta_{2}=0$ is not really a splitting method since it can be written as a one-step method

$$
\left\{\begin{array}{l}
r_{i}^{n+1}-r_{i}^{n}+a_{\star} \Delta t_{i+1}^{u_{i+1}^{n}-u_{i-1}^{n}}-\nu_{r} \Delta t^{\frac{r_{i+1}^{n}-2 r_{i}^{n}+r_{i-1}^{n}}{\Delta x^{2}}=0} \\
u_{i}^{n+1}-u_{i}^{n}+a_{\star} \Delta t_{i+1}^{n-r_{i-1}^{n}}-\nu_{u} \Delta t^{u_{i+1}^{n}-2 u_{i}^{n}+u_{i-1}^{n}} \\
\Delta x^{2}
\end{array}=\omega \Delta t\left[\theta_{1} v_{i}^{n}+\left(1-\theta_{1}\right) v_{i}^{n+1}\right], ~\left\{t u_{i}^{n+1} .\right.\right.
$$

## B Stability of the one step All-Froude Godunov scheme

Here we detail the proof of Lemma 6 .
Proof. The characteristic polynomial $\mathcal{P}(\lambda)$ of the amplification matrix (39) is given by

$$
\begin{aligned}
\mathcal{P}(\lambda)=\left(1-\lambda-2 \kappa_{r}\left|a_{\star}\right| \sigma s^{2}\right)\left(\lambda^{2}-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \lambda\right. & \left.+\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\right) \\
& +(1-\lambda) \frac{4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)},
\end{aligned}
$$

which can be decomposed as

$$
\begin{equation*}
\mathcal{P}(\lambda)=\mathcal{P}_{0}(\lambda)+\kappa_{r} \mathcal{P}_{1}(\lambda), \tag{45}
\end{equation*}
$$

where $\mathcal{P}_{0}$ is given by (41) and

$$
\mathcal{P}_{1}(\lambda)=-2\left|a_{\star}\right| \sigma s^{2}\left(\lambda^{2}-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \lambda+\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\right) .
$$

Since the roots of polynomial $\mathcal{P}_{0}$ are simple - see the proof of Lemma 5 - a classical continuity argument [14] allows us to write the roots of the polynomial $\mathcal{P}$ by using an asymptotic expansion

$$
\begin{equation*}
\lambda=\lambda^{(0)}+\kappa_{r} \lambda^{(1)}+\mathcal{O}\left(\kappa_{r}^{2}\right) \tag{46}
\end{equation*}
$$

where $\lambda^{(0)}$ is a root of $\mathcal{P}_{0}$. The stability of the scheme is obtained if the modulus of all roots of $\mathcal{P}$ is smaller than one. If $\lambda^{(0)}=\lambda_{ \pm}$, the results is obvious since one can ensure $\left|\lambda_{ \pm}\right|<1$ by considering

$$
\Delta t \leq K \min \left\{\Delta t_{a}, \Delta t_{b}\right\},
$$

with $K<1$ small enough and $\Delta t_{a}, \Delta t_{b}$ given in (36a-36b). The case $\lambda^{(0)}=\lambda_{0}=1$ is a bit more tricky. By inserting the asymptotic expansion (46) into relation (45), we obtain

$$
\mathcal{P}(\lambda)=\kappa_{r}\left[\lambda_{1} \mathcal{P}_{0}^{\prime}\left(\lambda_{0}\right)+\mathcal{P}_{1}\left(\lambda_{0}\right)\right]+\mathcal{O}\left(\kappa_{r}^{2}\right) .
$$

The condition $\mathcal{P}(\lambda)=0$ thus implies

$$
\lambda_{1}=-\frac{\mathcal{P}_{1}\left(\lambda_{0}\right)}{\mathcal{P}_{0}^{\prime}\left(\lambda_{0}\right)}
$$

Easy computations lead to

$$
\begin{aligned}
\mathcal{P}_{1}\left(\lambda_{0}\right) & =-2\left|a_{\star}\right| \sigma s^{2}\left(1-\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}+\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}\right) \\
& =-\frac{2\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \gamma^{2}<0 .
\end{aligned}
$$

On the other hand, since $\mathcal{P}_{0}(\lambda)=(1-\lambda) \widetilde{\mathcal{P}_{0}}(\lambda)$, we have

$$
\begin{aligned}
\mathcal{P}_{0}^{\prime}(1) & =-\widetilde{\mathcal{P}_{0}}(1)=-1+\frac{2-\gamma^{2}\left(\theta_{1}+\theta_{2}-2 \theta_{1} \theta_{2}\right)-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}-\frac{1+\gamma^{2} \theta_{1} \theta_{2}-2 \kappa_{u}\left|a_{\star}\right| \sigma s^{2}}{\Lambda\left(\theta_{1}, \theta_{2}\right)}-\frac{4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)} \\
& =-\frac{\gamma^{2}+4 a_{\star}^{2} \sigma^{2} s^{2}\left(1-s^{2}\right)}{\Lambda\left(\theta_{1}, \theta_{2}\right)}<0 .
\end{aligned}
$$

It follows that

$$
-2\left|a_{\star}\right| \sigma<\lambda_{1}<0
$$

and the scheme is stable.

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