ANALYSIS OF A HIGH-ORDER SPACE AND TIME DISCONTINUOUS GALERKIN METHOD FOR ELASTODYNAMIC EQUATIONS. APPLICATION TO 3D WAVE PROPAGATION.

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Abstract

In this paper, we introduce a high-order discontinuous Galerkin method, based on centered fluxes and a family of high-order leap-frog time schemes, for the solution of the 3D elastodynamic equations written in velocity-stress formulation. We prove that this explicit scheme is stable under a CFL type condition obtained from a discrete energy which is preserved in domains with free surface or decreasing in domains with absorbing boundary conditions. Moreover, we study the convergence of the method for both the semi-discrete and the fully discrete schemes, and we illustrate the convergence results by the propagation of an eigenmode. We also propose a series of absorbing conditions which allow improving the convergence of the global scheme. Finally, several numerical applications of wave propagation, using a 3D solver, help illustrating the various properties of the method.

AMS subject classifications: 35L50,35F10,35F15,35L05,35Q99

Key words: discontinuous Galerkin method, centered flux, leap-frog scheme, elastodynamic equation

Introduction

In the last few decades, the computational seismology has become an essential tool to simulate realistic wavefields of local, regional or even global wave propagation problems. The physics governing these phenomena is now well understood and many different accurate numerical methods have been developed and can deal with three-dimensional realistic applications thanks to a continuous increase of the computational resources and the use of parallel computational facilities.

Among all the numerical methods proposed for simulations in time domain, the most popular is undoubtedly the finite difference (FD) method and its many improvements from the initial FD schemes proposed by Alterman and Karal [7] or Kelly et al. [41] such as, for instance, the introduction of the velocity-stress system and the staggered-grids (Madariaga [47], Virieux [61]), fourth-order schemes in space (Bayliss et al. [9], Levander [45]) and rotated staggered-grids (Saenger et al. [57]) allowing strong fluctuations of the elastic parameters. If these techniques lead to simple and cheap algorithms, their major drawback is the restriction to cartesian grids not suited for geometrical internal or surfacic nonlinearities (topography) but, the use of non uniform grids (Pitarka [54]) or even discontinuous grids (Aoi and Fujiwara [6]) enabled improving the accuracy at the free surface and, more recently, an hybrid method coupling a finite-difference technique for the most part of the domain and a finite-element method in subdomains containing the nonlinearities (topography, faults) has been proposed (Moczo et al. [51], Galis et al. [33]). Some other methods have been further developed such as finite element (FE) methods which allow meshes adapted to complex geometries (Lysmer and Drake [46], Marfurt [48], Bao et al. [8]). However, they are often very costly because one needs to invert a global mass matrix at each time step which is not well adapted to the use of high-order space approximations. This difficulty was overcome by the use of Gauss-Lobato Legendre quadrature formulae at the root of the spectral element methods (SEM), usually applied
to second-order system in displacement (see, for instance Komatitsch and Vilotte [42], Chaljub et al. [14] amongst many contributions). However, except some applications to triangular meshes (Mercerat et al. [50], Mazzieri and Rappeti [49]) and resulting in non-diagonal mass matrices, these methods are mostly developed for quadrangular or hexahedral discretizations which may complicate the construction of meshes adapted to heterogeneous media. As an attractive alternative for a better representation of complex geometries, non conforming discontinuous Galerkin spectral element methods have been recently developed and analysed for the wave propagation in 2D quadrangular meshes (Antonietti et al. [5]) or 3D affine hexahedral meshes using a velocity-strain formulation (Wilcox et al. [63], Bui-Thanh and Ghattas [11]) for coupled elastic/acoustic media. When methods based on conforming meshes are used, simplicial meshes permits a better approximation of complex geometries (topography, faults). For this reason, we study a high-order Discontinuous Galerkin (DG) method applied to tetrahedral meshes.

The Discontinuous Galerkin (DG) method has been initially introduced by Reed and Hill [55] for the solution of neutron transport problems. A first analysis of the method for hyperbolic equations has been presented by Lesaint and Raviart [44]. Their result has been improved and broadened by Johnson et al. [36, 37]. The method is now very popular for solving many problems of physics governed by hyperbolic partial differential equations. Detailed reviews of methods and applications can be found in Cockburn et al. [22] or Hesthaven and Warburton [31].

This method can be seen as a discontinuous finite element method with a great flexibility in the choice of the local degree $p$ of the polynomial interpolation while being adapted to most meshes such as unstructured or non-conforming discretizations. Both aspects can be combined to build high-order hp-adaptive finite element methods (Süli et al. [59], Demkowicz and Kurtz [24]). The major advantage of this method is an easy extension to high-order in space realized by local basis functions defined in each element. The type of basis functions is not fixed and many families are available with their own advantages/drawbacks that can impact the accuracy (De Basabe et al. [23]). The inversion of a global mass matrix is avoided when an explicit scheme in time is used since all calculations are done at the local level producing block diagonal matrices. On the contrary, fluxes between neighboring elements must be accounted for and many strategies have been studied in the literature (for instance, see Hesthaven and Warburton [31]). More, it has a simple computational pattern with elements which makes the method naturally suited for parallelization.

Among the pioneering methods, Cockburn and Shu [18] introduced a class of high-order discontinuous Galerkin finite element methods, based on upwind fluxes and a Runge-Kutta time scheme, for 1D scalar conservation laws. The method was then extended to systems and multidimensional cases in [19, 20, 21] and the analysis proved the good convergence properties of the method. Other properties are also important to define an accurate scheme. The study of the dispersive and dissipative properties of a method provides information on its ability to accurately propagate a wave, especially high frequency, along many wavelengths (as it is the case for simulation of realistic problems). Hu et al. [32] studied dispersive and dissipative properties of a DG method for both 1D and 2D equations. They concluded that they depend on the type of fluxes: the error being dominated by dissipation for upwind fluxes whereas the dissipation rate is zero for centered fluxes. Sherwin [58] compared the dispersive properties of both continuous and discontinuous formulations. Ainsworth [3] and Ainsworth et al. [4] analyzed these properties for high-order DG finite element methods, considering linear advection and second-order wave equations respectively. One of the results of these studies is that it may be advantageous to increase the order of the method instead of refining the mesh. An important issue is then to study how to choose the order of polynomial basis function for each element in the computational mesh to obtain a predetermined error. Lähivaara and Huttunen [43] studied such a criterion, connecting the polynomial order $p$ to the mesh parameter $h$, for acoustic and elastic wave equation. Chevaugeon et al. [16] proposed a number of rules to choose an optimal order for the spatial approximation and an optimal time-stepping scheme in order to obtain accurate solutions of the linearized Euler equations in complex geometries and then reduce the size of the problem and the CPU time.

An analysis of the spatial discretization error of a DG method associated with the solutions of 1D hyperbolic conservation laws has been proposed by Adjerid et al. [1]. In particular, superconvergence results of the DG method at Radau points is obtained. These results are completed by Cheng and Shu [15] who studied the evolution of the space error with time and Zhong and Shu [64] who considered the fully discretized approximation.

In spite of its success in many domains of applications, the DG method has been applied relatively late to seismic wave propagation problems in comparison to other domains of application. If we restrict ourselves to methods designed for the velocity-stress first-order formulation, Käser and Dumbser [38] presented the ADER-DG method for 2D applications that has been widely developed and extended to 3D configurations (Dumbser and Käser [26]), viscoelastic media (Käser et al. [39]), applied to a quantitative accuracy analysis (Käser et al. [40]) or dynamic rupture simulations (Pelties et al. [52]), to mention only a few of the many developments related to this method. The ADER-DG method is based on upwind fluxes and the Cauchy-Kowalewski procedure that
automatically provides schemes with the same space and time level accuracy which allow achieving very high levels of convergence. Other references are related to a non-diffusive DG scheme based on centered fluxes and a leap-frog time scheme: Delcourte et al. [25] for the introduction and some properties of the method, Etienne et al. [28] for 3D simulations in a realistic basin, Tago et al. [60] for 3D rupture simulations or Peyrusse et al. [53] for 2D viscoelastic media.

In this paper, we study the P-SV seismic wave propagation considering an isotropic, linearly elastic medium by solving the velocity-stress formulation of the elastodynamic equations. For the discretization of this system, we apply the method described in Delcourte et al. [25], based on centered fluxes and a leap-frog time-discretization. The extension to high-order in space is realized by Lagrange polynomial functions (of degree 0 to 4). Starting with the conclusions of the previous paper [25], we wish to improve some particular points. A numerical convergence study proved that the accuracy of the time scheme is crucial when global high-order accuracy is sought. Indeed, unlike the ADER method that has the same level of accuracy in space and time, the standard leap-frog time scheme, second-order accurate, reduces the global convergence order limiting the interest for high-order space discretizations. Then, we propose an extension of the leap-frog scheme to higher (even) orders of accuracy, following a method proposed for the Maxwell equations by Young [62] and applied in the DG context by Fahs [29]. This method allows us to achieve temporal accuracy to any even order desired, when free surface conditions are applied, by introducing an iterative procedure. Another technique, based on the modified equation approach, has been applied by Agut et al. [2] to construct high-order time schemes for the second-order wave equation. If our analysis is done for a general class of leap-frog schemes, we restrict the numerical application to standard and fourth-order cases, since fourth degree polynomial functions have been considered at the most for spatial discretization. This complete a first introduction of the method, in a two-dimensional context, done by Glinsky et al. [35], especially by adding an analysis of the scheme.

A second limitation of the method found in [25] concerns the low-order absorbing boundary conditions. Indeed, when considering the error evolution during the propagation of a wave in an artificially bounded domain, we can notice a reduction in the convergence due to the absorbing boundary condition. The convergence studies are generally done for problems involving free surface conditions or periodicity conditions, as it is the case for the numerous convergence studies using the ADER method. Our objective here is not to study this point in detail; as mentionned in the review done by Givoli [34], this complex topic is a study in itself. But some simple improvements can be obtained by the evaluation of the absorbing boundary terms at some particular time levels and we wish to study this particular point.

This article is organized as follows. In section 1, we state the velocity-stress formulation in a symmetric pseudo-conservative form. Then, in section 2, we detail the discretization of the equations system by the Discontinuous Galerkin method. The high-order leap-frog time scheme as well as the approximation of the boundary conditions are also presented. After, we study, in section 3, the preservation of a discrete elastodynamic energy and the stability of the scheme, taking into account the free surface or the absorbing boundary conditions. Section 4 is devoted to the convergence analysis of the semi-discrete and fully discretized schemes. Finally, we illustrate this study in section 5 by some basic numerical applications in three dimensions of space, with the aim of validating the 3D solver and studying, numerically, particular points of the method: the propagation of an eigenmode for a convergence study and a comparison between standard and fourth-order leap-frog time-schemes, the propagation of a pulse for a comparison of two different absorbing boundary conditions and the propagation of an explosive source in a half-space. The main objective of this paper being the analysis of the method, applications to 3D more complex configurations, especially in heterogeneous media requiring adapted unstructured meshes, will be the object of a further publication.

1 Velocity-stress formulation in pseudo-conservative form

The P-SV wave propagation in an isotropic, linearly elastic medium is modelised by the elastodynamic equations, which initially write in displacement-stress formulation; let be $\vec{U} = (U_\alpha)_{\alpha \in \{x,y,z\}}$ the displacement vector and $\bar{\sigma} = (\sigma_{\alpha,\beta})_{\alpha,\beta \in \{x,y,z\}}$, the stress tensor, then the system reads

$$\begin{cases} \rho \frac{\partial^2 \vec{U}}{\partial t^2} = \nabla \cdot \bar{\sigma}, \\ \bar{\sigma} = \lambda (\nabla \cdot \vec{U}) \vec{T} + \mu \left( \nabla \vec{U} + (\nabla \vec{U})^t \right), \end{cases} \tag{1}$$
where $\mathbf{I}$ is the identity matrix, $\rho$ is the density of the medium and $\lambda$ and $\mu$ are the Lamé parameters related to shear and compressional velocities ($v_s$ and $v_p$) in the medium by $v_s = \sqrt{\frac{\mu}{\rho}}$ and $v_p = \sqrt{\frac{\lambda + 2\mu}{\rho}}$.

We introduce the velocity vector $\mathbf{V} = \partial_t \mathbf{U}$ in the equation (1) and obtain the first-order velocity-stress formulation [61]:

$$
\begin{align*}
\rho \partial_t \mathbf{V} &= \nabla \cdot \mathbf{\sigma}, \\
\partial_t \mathbf{\sigma} &= \lambda (\nabla \cdot \mathbf{V}) \mathbf{I} + \mu (\nabla \mathbf{V} + (\nabla \mathbf{V})^T).
\end{align*}
$$

Let $\mathbf{W} = (\mathbf{V}, \sigma)^T$ be the vector composed of the velocity components $\mathbf{V} = (V_x, V_y, V_z)^T$ and the stress components $\mathbf{\sigma} = (\sigma_{xx}, \sigma_{yy}, \sigma_{zz}, \sigma_{xy}, \sigma_{xz}, \sigma_{yz})^T$, since the stress tensor is symmetric, then, the system (2) can be rewritten as

$$
\partial_t \mathbf{W} + \sum_{\alpha \in \{x,y,z\}} A_{\alpha} (\rho, \lambda, \mu) \partial_{\alpha} \mathbf{W} = 0.
$$

We choose here not to detail the matrices $A_{\alpha} (\rho, \lambda, \mu)$.

In order to express the system (3) in a pseudo-conservative form, we introduce the following change of variables on the stress components,

$$
\tilde{\mathbf{\sigma}} = \mathbf{T} \mathbf{\sigma} = \left(\frac{\sigma_{xx} + \sigma_{yy} + \sigma_{zz}}{3}, \frac{2\sigma_{xx} - \sigma_{yy} - \sigma_{zz}}{3}, -\sigma_{xx} + 2\sigma_{yy} - \sigma_{zz}, \sigma_{xy}, \sigma_{xz}, \sigma_{yz}\right).
$$

So, the system reads

$$
\begin{align*}
\rho \partial_t \mathbf{V} &= \sum_{\alpha \in \{x,y,z\}} M_{\alpha} \partial_{\alpha} \mathbf{\sigma}, \\
A_0 (\lambda, \mu) \partial_t \mathbf{\sigma} &= \sum_{\alpha \in \{x,y,z\}} N_{\alpha} \partial_{\alpha} \mathbf{\sigma},
\end{align*}
$$

where $A_0 (\lambda, \mu) = \text{diag} \left(\frac{3}{3\lambda + 2\mu}, \frac{3}{2\mu}, \frac{3}{\mu}, 1, 1, 1\right)$ is a diagonal matrix containing the characteristics of the medium and

$$
\begin{align*}
M_x &= \begin{pmatrix} 1 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 \end{pmatrix}, & M_y &= \begin{pmatrix} 0 & 0 & 1 & 0 & 0 & 0 \\ 1 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 \end{pmatrix}, & M_z &= \begin{pmatrix} 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 \\ 1 & -1 & 0 & 0 & 0 & 0 \end{pmatrix}, \\
N_x &= \begin{pmatrix} 1 & 0 & 0 \\ 2 & 0 & 0 \\ -1 & 0 & 0 \end{pmatrix}, & N_y &= \begin{pmatrix} 0 & 1 & 0 \\ 0 & -1 & 0 \\ 0 & 2 & 0 \end{pmatrix}, & N_z &= \begin{pmatrix} 0 & 0 & 1 \\ 0 & 0 & 1 \\ 0 & 0 & 1 \end{pmatrix}.
\end{align*}
$$

We can notice that now the matrices $M_\alpha$ and $N_\alpha$ ($\alpha = x, y, z$) are constant and do not depend anymore on the material properties. So, the system (5) is a pseudo-conservative formulation of (3). At last, we multiply the second equation of (5) by the following positive definite symmetric (PDS) matrix

$$
S = \begin{pmatrix} 1 & 0 & 0 & 0 & 0 & 0 \\ 0 & \frac{3}{5} & \frac{1}{5} & 0 & 0 & 0 \\ 0 & 0 & \frac{3}{5} & \frac{2}{5} & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 & 1 \end{pmatrix}
$$

in order to obtain a symmetric system. Therefore, we finally get the symmetric pseudo-conservative formulation:

$$
\begin{align*}
\rho \partial_t \mathbf{V} &= \sum_{\alpha \in \{x,y,z\}} M_{\alpha} \partial_{\alpha} \mathbf{\sigma}, \\
A (\lambda, \mu) \partial_t \mathbf{\sigma} &= \sum_{\alpha \in \{x,y,z\}} M_{\alpha}^T \partial_{\alpha} \mathbf{\sigma},
\end{align*}
$$

where $\mathbf{\sigma}$ is the stress tensor in the symmetric pseudo-conservative formulation.
where $\Lambda(\lambda, \mu) = S \Lambda_0(\lambda, \mu)$ is a positive definite symmetric matrix whose spectrum is

$$Sp(\Lambda) = \left\{ \frac{1}{2\mu}, \frac{3}{2\mu}, \frac{1}{\mu}, \frac{1}{\mu}, \frac{3}{3\lambda + 2\mu} \right\} \subset \mathbb{R}^+_*$$

and we can notice that $\mathbb{S}N_\alpha = M^t\alpha (\alpha = x, y, z)$. This formulation will be very useful to establish the energy preservation. Then, we add:

(i) a physical boundary condition on the free surface

$$\bar{\sigma} \mathbf{n} = 0,$$  

which will be rewritten by respect to $\bar{\sigma}$ in section 2.3.

(ii) absorbing boundary conditions

$$\mathbb{P} \bar{\sigma} = -A \bar{V} \text{ and } \mathbb{P}^t \bar{V} = -B \bar{\sigma},$$

(where the matrices $\mathbb{P}$, $A$ and $B$ will be specified in section 2.3) to approximate an infinite domain.

At last, we need to impose initial data

$$\bar{V}(0) = \bar{V}_0 \text{ and } \bar{\sigma}(0) = \bar{\sigma}_0,$$  

(9)

to ensure the existence and the uniqueness of the solution $(\bar{V}, \bar{\sigma})$ of the system (6)-(7)-(8)-(9).

2 A discontinuous Galerkin method combined with a leap-frog scheme

2.1 Integration on a simplex

We consider a bounded polyhedral domain $\Omega$ of $\mathbb{R}^3$, discretized in $N_T$ tetrahedra $T_i$, which form a partition of the domain. We assume that the characteristics of the medium are constant over each element $T_i$ and denoted by $(\rho_i, \lambda_i, \mu_i)$. We multiply, in the sense of the scalar product, the first (resp. the second) equation of the problem (6) by a vector field $\varphi_i \in \mathbb{R}^3$ (resp. $\psi_i \in \mathbb{R}^6$) and we integrate them on each element $T_i$. Then, we apply the Green formula to the right-hand sides:

$$\rho_i \int_{T_i} \varphi_i \partial_t \bar{V} \, dv = \sum_{\alpha \in \{x,y,z\}} \left[ - \int_{\partial T_i} (\partial_n \varphi_i)^\alpha M_\alpha \bar{\sigma} \, ds + \int_{\partial T_i} \varphi_i^\alpha M_\alpha \partial_t n_\alpha \, ds \right],$$

(10a)

$$\int_{T_i} \bar{\psi}_i \Lambda_i \partial_t \bar{\sigma} \, dv = \sum_{\alpha \in \{x,y,z\}} \left[ - \int_{\partial T_i} (\partial_n \bar{\psi}_i)^\alpha M_\alpha \bar{V} \, ds + \int_{\partial T_i} \bar{\psi}_i^\alpha M_\alpha \partial_t \bar{V} \, ds \right],$$

(10b)

where we set $\Lambda_i = \Lambda(\lambda_i, \mu_i)$ and $\mathbf{n}_i = (n_x, n_y, n_z)^t$ represents the outwards unit normal vector to $T_i$.

2.2 Evaluation of volume integrals

The approximations of $\bar{V}$ and $\bar{\sigma}$ are denoted by the fields $\bar{V}_h$ and $\bar{\sigma}_h$ which are defined locally on each element $T_i$ and may be discontinuous through the interfaces, so that we set

$$\forall i, \bar{V}_{h|T_i} = \bar{V}_i \text{ and } \bar{\sigma}_{h|T_i} = \bar{\sigma}_i.$$  

(11)

Let us denote $\mathcal{P}_m(T_i)$ the set of polynomials over $T_i$ with a degree $m$ and consider the vector fields $\varphi_{ij} \in \mathcal{P}_m^3(T_i)$ and $\psi_{ij} \in \mathcal{P}_m^6(T_i)$ whose each component is a Lagrange nodal interpolant, then we can write $\bar{V}_i$ and $\bar{\sigma}_i$ as linear combinations of time-dependent fields:

$$\bar{V}_i(x, y, z, t) = \sum_{j=1}^{dof} \hat{\varphi}_{ij}(t) \varphi_{ij}(x, y, z) \text{ and } \bar{\sigma}_i(x, y, z, t) = \sum_{j=1}^{dof} \hat{\psi}_{ij}(t) \psi_{ij}(x, y, z),$$

(12)

where $dof$ denotes the number of degrees of freedom on the element $T_i$. To approximate the volume integrals on $T_i$, we just have to replace the fields $\bar{V}$ and $\bar{\sigma}$ in (10) by $\bar{V}_h$ and $\bar{\sigma}_h$. 

5
2.3 Approximation on faces

To calculate the integrals on $\partial \mathcal{T}_i$ of (10), we split this boundary into internal and boundary faces. We denote by $\mathcal{V}(i)$ the set of indices of the neighboring elements of $\mathcal{T}_i$ and we note by $S_{ik}$ each internal face common to both elements $\mathcal{T}_i$ and $\mathcal{T}_k$, i.e. $S_{ik} = \mathcal{T}_i \cap \mathcal{T}_k$. Moreover, some elements $\mathcal{T}_i$ have one or more faces common to the boundary of the domain. The set of the indices $k$ of such faces $S_{ik} = \mathcal{T}_i \cap \partial \Omega$ is denoted by $\mathcal{K}(i)$ for absorbing boundaries and $\mathcal{K}(i)$ for free surface boundaries. Remark that, for most elements, $\mathcal{E}(i)$ and $\mathcal{K}(i)$ are empty sets. Therefore, introducing $P_{ik} = \sum_{\alpha \in \{x,y,z\}} M_{\alpha} n_{\alpha ik}$, we have:

\[
\sum_{\alpha \in \{x,y,z\}} \int_{\partial \mathcal{T}_i} \varphi^\alpha M_{\alpha} \bar{\sigma} n_{\alpha} ds = \sum_{k \in \mathcal{V}(i)} \int_{S_{ik}} \varphi^\alpha P_{ik} \bar{\sigma} ds + \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \varphi^\alpha P_{ik} \bar{\sigma} ds + \sum_{k \in \mathcal{K}(i)} \int_{S_{ik}} \varphi^\alpha P_{ik} \bar{\sigma} ds, \tag{13a}
\]
\[
\sum_{\alpha \in \{x,y,z\}} \int_{\partial \mathcal{T}_i} \psi^\prime^\alpha M_{\alpha} \bar{v} n_{\alpha} ds = \sum_{k \in \mathcal{V}(i)} \int_{S_{ik}} \psi^\prime^\alpha P_{ik} \bar{v} ds + \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \psi^\prime^\alpha P_{ik} \bar{v} ds + \sum_{k \in \mathcal{K}(i)} \int_{S_{ik}} \psi^\prime^\alpha P_{ik} \bar{v} ds, \tag{13b}
\]

where $\bar{n}_{ik} = (n_{x ik}, n_{y ik}, n_{z ik})^t$ represents the unit normal vector of $S_{ik}$, oriented from $\mathcal{T}_i$ towards $\mathcal{T}_k$.

For the **interior faces** $S_{ik}$ of $\mathcal{T}_i (k \in \mathcal{V}(i))$, we choose to apply centered fluxes by introducing the mean-value on this face:

\[
\bar{V}_{S_{ik}} \simeq \frac{1}{2} \left( \bar{V}_i + \bar{V}_k \right) \quad \text{and} \quad \bar{\sigma}_{S_{ik}} \simeq \frac{1}{2} \left( \bar{\sigma}_i + \bar{\sigma}_k \right). \tag{14}
\]

For the boundary integrals, two types of boundary conditions have been considered: a free surface condition at the physical interface between air and the medium, and an absorbing condition on the artificial boundaries of an infinite domain.

**Free surfaces:** On these faces, we compute the fluxes by introducing weakly the physical condition (7) in the third term of (13a) and the physical condition is rewritten in the new variable $\bar{\sigma}$ via the change of variables (4). No physical condition is applied to the velocity in the third term of (13b). So, for a boundary face $S_{ik}$ of $\mathcal{T}_i$, with $k \in \mathcal{K}(i)$, this condition reduces to:

\[
\bar{V}_{S_{ik}} \simeq \bar{V}_i \quad \text{and} \quad P_{ik} \bar{\sigma}_{S_{ik}} = 0. \tag{15}
\]

**Absorbing surfaces:** To simulate infinite domains, we introduce artificial boundaries, on which we impose absorbing conditions. Therefore, for any $k \in \mathcal{E}(i)$ and any real unit vector $\vec{n}_{ik}$, we define the matrix

\[
A_{n_{ik}}(\rho_i, \lambda_i, \mu_i) = \sum_{\alpha \in \{x,y,z\}} A_{\alpha}(\rho_i, \lambda_i, \mu_i) n_{\alpha ik} = -\left( \begin{array}{ccc} 0_{3 \times 3} & \frac{1}{\rho_i} & \frac{1}{\rho_i} P_{ik} \\ \Lambda_i^{-1} P_{ik} & 0_{6 \times 6} & \end{array} \right), \tag{16}
\]

which is diagonalizable in $\mathbb{R}$, i.e. all its eigenvalues $e_k (k=1,...,9)$ are real:

\[
e_1 = -v_{p_i}, \quad e_2 = e_3 = -v_{s_i}, \quad e_4 = e_5 = e_6 = 0, \quad e_7 = e_8 = v_{s_i}, \quad e_9 = v_{p_i},
\]

with $v_{p_i} = v_p(\rho_i, \lambda_i, \mu_i) = \sqrt{\lambda_i + 2 \mu_i / \rho_i}$ and $v_{s_i} = v_s(\rho_i, \lambda_i, \mu_i) = \sqrt{\mu_i / \rho_i}$,

and we note by $P_{\alpha}(\rho, \lambda, \mu)$ the matrix whose column $k$ is the right eigenvector associated to the eigenvalue $e_k$:

\[
\begin{pmatrix}
3\sqrt{\nu_x} & -v_{x} n_{x} n_{y} & v_{x} n_{01, -1} & 0 & 0 & 0 & -v_{x} n_{00, -1} & v_{x} n_{y} & -3 \sqrt{\nu_x} \\
3\sqrt{\nu_y} & v_{y} n_{01, -1} & -v_{y} n_{x} n_{y} & 0 & 0 & 0 & v_{y} n_{x} n_{y} & -v_{y} n_{y} & -3 \sqrt{\nu_y} \\
3\sqrt{\nu_z} & v_{z} n_{x} n_{z} & 0 & 0 & 0 & 0 & -v_{z} n_{y} n_{z} & v_{z} n_{y} & -3 \sqrt{\nu_z} \\
3\lambda + 2 \mu & 0 & -n_{x} n_{00, -1} & n_{y} n_{01, -1} & -n_{x} n_{11, 0} & -n_{x} n_{11, 0} & 0 & 0 & 3\lambda + 2 \mu \\
2\mu n_{x} n_{y} & 2 \mu n_{00, -1} & -2 \mu n_{x} n_{y} & 0 & 0 & 0 & 0 & 0 & -2 \mu n_{y} n_{x} \\
-2 \mu n_{x} n_{y} & 2 \mu n_{01, -1} & 0 & 0 & 0 & -2 \mu n_{x} n_{y} & 0 & 0 & -2 \mu n_{y} n_{x} \\
6 \mu n_{y} n_{y} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 6 \mu n_{y} n_{y} \\
6 \mu n_{y} n_{z} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 6 \mu n_{y} n_{z} \\
6 \mu n_{y} n_{z} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 6 \mu n_{y} n_{z}
\end{pmatrix}
\]
where, to simplify, we set \( n_{a,b,c} = an_x^2 + bn_y^2 + cn_z^2 \). Then, the boundary absorbing conditions consist in an upwind technique where we only take into account the outgoing waves neglecting the ingoing waves part. For that, we approximate \( A_{n,ik} \) by the matrix \( A_{n,ik}^+ \) where we set \( A_{n,ik}^+ = (P_{n,ik}D^TP_{n,ik}^{-1}) (\rho_i, \lambda_i, \mu_i) \) with \( D^+ (\rho_i, \lambda_i, \mu_i) \) the diagonal matrix composed by the positive eigenvalues of \( A_{n,ik} \). Therefore, according to (16), we have to approximate

\[
\begin{pmatrix}
\bar{P}_{ik} \vec{\tau} \\
\bar{P}_{ik} \vec{V}
\end{pmatrix} = -\begin{pmatrix}
\rho_i \bar{T}_3 \\
\rho_i \bar{T}_3
\end{pmatrix} \begin{pmatrix}
0_{3\times 6} \\
0_{3\times 6}
\end{pmatrix} \Lambda_i \begin{pmatrix}
A_{n,ik} (\rho_i, \lambda_i, \mu_i) \\
A_{n,ik}^+ (\rho_i, \lambda_i, \mu_i)
\end{pmatrix} \begin{pmatrix}
\vec{V} \\
\vec{V}_i
\end{pmatrix},
\]

where \( \bar{T}_3 \) denotes the identity matrix, by

\[
\begin{pmatrix}
\bar{P}_{ik} \vec{\tau} \\
\bar{P}_{ik} \vec{V}
\end{pmatrix} \simeq -\begin{pmatrix}
\rho_i \bar{T}_3 \\
\rho_i \bar{T}_3
\end{pmatrix} \begin{pmatrix}
0_{3\times 6} \\
0_{3\times 6}
\end{pmatrix} \Lambda_i \begin{pmatrix}
A_{n,ik}^+ (\rho_i, \lambda_i, \mu_i) \\
A_{n,ik} (\rho_i, \lambda_i, \mu_i)
\end{pmatrix} \begin{pmatrix}
\vec{V} \\
\vec{V}_i
\end{pmatrix}.
\]

After calculus, we obtain:

\[
\begin{pmatrix}
\rho_i \bar{T}_3 \\
0_{3\times 6}
\end{pmatrix} \begin{pmatrix}
A_{n,ik}^+ (\rho_i, \lambda_i, \mu_i) \\
A_{n,ik} (\rho_i, \lambda_i, \mu_i)
\end{pmatrix} = \frac{1}{2} \begin{pmatrix}
A_{ik} & -R_{ik} \\
-\bar{R}_{ik} & \bar{P}_{ik}
\end{pmatrix}
\]

where

\[
\begin{pmatrix}
A_{ik} \\
\bar{R}_{ik}
\end{pmatrix} = \begin{pmatrix}
\frac{(v_{pi} - v_{si})n_{x,ik}^2 + v_{si}}{(v_{pi} - v_{si})n_{x,ik}n_{y,ik}} \\
\frac{(v_{pi} - v_{si})n_{x,ik}n_{y,ik} + v_{si}}{(v_{pi} - v_{si})n_{y,ik}n_{z,ik}}
\end{pmatrix}
\]

and

\[
\begin{pmatrix}
\bar{R}_{ik} \\
\bar{P}_{ik}
\end{pmatrix} = \begin{pmatrix}
n_{x,ik}^2 - n_{z,ik}^2 & n_{y,ik}^2 - n_{z,ik}^2 & 2n_{x,ik}n_{y,ik} & 2n_{x,ik}n_{z,ik} & 2n_{y,ik}n_{z,ik} \\
n_{y,ik}^2 - n_{z,ik}^2 & n_{x,ik}^2 - n_{z,ik}^2 & n_{x,ik}n_{y,ik} & n_{x,ik}n_{z,ik} & n_{y,ik}n_{z,ik} \\
n_{x,ik}^2 - n_{z,ik}^2 & n_{y,ik}^2 - n_{z,ik}^2 & n_{x,ik}n_{y,ik} & n_{x,ik}n_{z,ik} & n_{y,ik}n_{z,ik} \\
2n_{x,ik}n_{y,ik} & n_{x,ik}n_{y,ik} & n_{x,ik}n_{y,ik} & n_{x,ik}n_{y,ik} & n_{y,ik}n_{z,ik} \\
2n_{y,ik}n_{z,ik} & n_{y,ik}n_{z,ik} & n_{y,ik}n_{z,ik} & n_{y,ik}n_{z,ik} & n_{y,ik}n_{z,ik}
\end{pmatrix}
\]

Finally, for a boundary face \( S_{ik} \) of \( \Gamma_i \) with \( k \in \mathcal{E}(i) \), the fluxes are approximated in (13) by:

\[
\int_{S_{ik}} \vec{\varphi} \bar{P}_{ik} \vec{\tau} \, ds \simeq \int_{S_{ik}} \frac{\bar{\varphi}_i}{2} \left[ -A_{ik} \vec{V}_i + \bar{P}_{ik} \vec{\tau}_i \right] \, ds,
\]

\[
\int_{S_{ik}} \bar{\varphi}_i \bar{P}_{ik} \vec{V}_i \, ds \simeq \int_{S_{ik}} \frac{\bar{\varphi}_i}{2} \left[ \bar{P}_{ik} \vec{V}_i - \bar{B}_{ik} \vec{\tau}_i \right] \, ds.
\]

**Remark 1** The absorbing boundary conditions (8) follow from (19)-(20). It is a first-order approximation, efficient for waves with a normal incidence to the artificial boundaries but whose accuracy could be improved especially to limit reflections in presence of grazing waves.

**Remark 2** \( A_{ik} \) is a symmetric positive definite matrix whose spectrum is

\[
Sp(A_{ik}) = \{ \rho_i v_{pi}, \rho_i v_{si}, \rho_i v_{si} \},
\]

whereas \( B_{ik} \) is a symmetric semi-definite positive matrix. Indeed, we have

\[
\det(B_{ik} - xI_b) = x^3(x^3 + ax^2 + bx + c)
\]
where

\[ a = -\frac{1}{\rho_k v_p v_s} \left[ v_p (2 - 2n^2_{yik} + n^2_{zik} - n^4_{yik} + 2n^4_{yik} + 2n^2_{yik}n^2_{zik}) + v_s (2 + n^4_{yik} - 2n^4_{yik} - 2n^2_{yik}n^2_{zik}) \right], \]

\[ b = \frac{1}{\rho_k^2 v_p v_s^2} \left[ v_p (1 + 2n^4_{yik} - 2n^2_{yik} - n^4_{zik} + n^2_{yik} - n^2_{yik}n^2_{zik} + 3n^2_{yik}n^4_{zik} + 3n^4_{yik}n^2_{zik}) 
+ v_s (4 - n^2_{yik} + 2n^4_{yik} + 4n^2_{yik}n^2_{zik} - 3n^2_{yik}n^4_{zik} - 3n^4_{yik}n^2_{zik} + 3n^4_{yik}n^2_{zik}) \right], \]

\[ c = -\frac{1}{\rho_k^3 v_p v_s^3} (2 + 2n^2_{zik} - 3n^4_{yik} + 3n^4_{yik} - n^4_{zik} - n^2_{yik}n^2_{zik} + 4n^4_{yik}n^2_{zik} + 4n^2_{yik}n^4_{zik}). \]

We can deduce that 0 is a third order eigenvalue of \( \mathbb{E}_{ik} \), but the three other eigenvalues are not easily specified. Then, we shall use the Sylvester’s criteria. For that, we denote by \( \mathbb{E}_{ik} = (b_{ij})_{1 \leq i, j \leq 6} \) the coefficients of the matrix \( \mathbb{E}_{ik} \) and \( \mathbb{B}_q = (b_{ji})_{1 \leq j, i \leq q} \) the submatrices of \( \mathbb{E}_{ik} \), and we observe that

\[ \det(\mathbb{E}_1) = \frac{1}{\rho_k^3 v_p v_s^3} > 0, \quad \det(\mathbb{E}_2) = \frac{1}{\rho_k^7 v_p v_s^7} \left[ 4n^2_{zik} n^2_{zik} + n^2_{zik} (n^2_{zik} + n^2_{zik}) \right] \geq 0, \]

\[ \det(\mathbb{E}_3) = \frac{9}{\rho_k^7 v_p v_s^7} n^2_{zik} n^4_{yik} n^2_{zik} \geq 0, \quad \det(\mathbb{E}_4) = \det(\mathbb{E}_5) = \det(\mathbb{E}_6) = 0, \]

so that the three other eigenvalues of \( \mathbb{E}_{ik} \) are necessarily positive.

### 2.4 Time-discretization

This subsection is devoted to the high-order leap-frog time-integration of the scheme in the same spirit as [29] for the Maxwell equations.

#### 2.4.1 The semi-discretized scheme

Taking into account of (11)-(14)-(15)-(23)-(24) in the equations (10a)-(10b), we obtain the following semi-discretized scheme for all \( T_i \), with \( i = 1, \ldots, N_T \):

\[ \rho_i \int_{T_i} \tilde{\phi}_i \partial_t \tilde{V}_i dv = - \sum_{\alpha \in \{x,y,z\}} \int_{T_i} (\partial_\alpha \tilde{\phi}_i)^t M_{\alpha i} \tilde{\sigma}_i dv + \sum_{k \in V_{(i)}} \int_{S_{ik}} \tilde{\phi}_i^t P_{ik} \tilde{\sigma}_i^t + \tilde{\sigma}_k^t \frac{1}{2} ds + \sum_{k \in \Xi_{(i)}} \int_{S_{ik}} \tilde{\phi}_i^t \left[ \frac{1}{2} P_{ik} \tilde{\sigma}_i - \frac{1}{2} \kappa_{ik} \tilde{V}_i \right] ds, \]  
\[ \int_{T_i} \tilde{\psi}_i^t \Lambda_i \partial_t \tilde{\sigma}_i dv = - \sum_{\alpha \in \{x,y,z\}} \int_{T_i} (\partial_\alpha \tilde{\psi}_i)^t M_{\alpha i} \tilde{V}_i dv + \sum_{k \in V_{(i)}} \int_{S_{ik}} \tilde{\psi}_i^t P_{ik} \tilde{V}_i + \tilde{V}_k \frac{1}{2} ds + \sum_{k \in \Xi_{(i)}} \int_{S_{ik}} \tilde{\psi}_i^t \left[ \frac{1}{2} P_{ik} \tilde{V}_i - \frac{1}{2} \mu_{ik} \tilde{\sigma}_i \right] ds + \sum_{k \in \Xi_{(i)}} \int_{S_{ik}} \tilde{\psi}_i^t P_{ik} \tilde{V}_i ds, \]  

According to (12), for each tetrahedron \( T_i \), we denote by \( \tilde{\psi}_i \) and \( \tilde{\sigma}_i \) the column vectors \( (\tilde{V}_{ij})_{1 \leq j \leq \text{dof}} \) and \( (\tilde{\sigma}_{ij})_{1 \leq j \leq \text{dof}} \) respectively. Replacing \( \phi_i \) by \( \phi_i^t \) in (25) and \( \psi_i \) by \( \psi_i^t \) in (26) for \( 1 \leq j \leq \text{dof} \), we obtain the following matricial form:

\[ M_{\phi i} \partial_t \tilde{V}_i = - \sum_{\alpha \in \{x,y,z\}} J_{\alpha i} \tilde{\sigma}_i + \sum_{k \in V_{(i)}} (K_{ik} \tilde{\sigma}_i + G_{ki} \tilde{\sigma}_k) + \sum_{k \in \Xi_{(i)}} (K_{ik} \tilde{\sigma}_i - 2 \alpha_{ik} \tilde{V}_i), \]

\[ M_{\Lambda i} \partial_t \tilde{\sigma}_i = - \sum_{\alpha \in \{x,y,z\}} L_{\alpha i} \tilde{V}_i + \sum_{k \in V_{(i)}} (K_{ik}^t \tilde{V}_i - G_{ki}^t \tilde{\sigma}_k) + \sum_{k \in \Xi_{(i)}} (K_{ik}^t \tilde{V}_i - 2 \beta_{ik} \tilde{\sigma}_i) + 2 \sum_{k \in \Xi_{(i)}} K_{ik}^t \tilde{V}_i, \]
where for all $1 \leq j, l \leq \text{dof}$, we have set

\[
(M_{\rho})_{jl} = \rho_i \int_{T_i} \varphi_{ij} \varphi_{il} \, ds, \quad (M_{\Lambda})_{jl} = \int_{T_i} \varphi_{ij} \Lambda_j \varphi_{il} \, ds, \\
(J_{\rho})_{jl} = \int_{T_i} (\partial_{\alpha} \varphi_{ij})^\alpha M_{\alpha} \varphi_{il} \, ds, \quad (L_{\rho})_{jl} = \int_{T_i} (\partial_{\alpha} \varphi_{ij})^\alpha M_{\alpha} \varphi_{il} \, ds, \\
(K_{ik})_{jl} = \frac{1}{2} \int_{S_{ik}} \varphi_{ij}^\alpha P_{ik} \varphi_{il} \, ds, \quad (G_{ik})_{jl} = \frac{1}{2} \int_{S_{ik}} \varphi_{ij}^\alpha P_{ik} \varphi_{kl} \, ds, \\
(\sigma_{ik})_{jl} = \frac{1}{2} \int_{S_{ik}} \varphi_{ij} \sigma_{ik} \varphi_{il} \, ds, \quad (R_{ik})_{jl} = \frac{1}{2} \int_{S_{ik}} \varphi_{ij} \sigma_{ik} \varphi_{il} \, ds,
\]

and used the equalities $\varphi_{ij}^\alpha P_{ik} \varphi_{il} = \varphi_{ij}^\alpha P_{ik} \varphi_{kl} = -\varphi_{ij}^\alpha P_{ki} \varphi_{jl}$, thanks to the orientation from $T_i$ to $T_k$ of $\nu_{\alpha ij}$, $\alpha \in \{x, y, z\}$, contained in $P_{ik}$. On the other hand, thanks to the Green formula, we have the following matricial relation for all $i \in [1, N_T]$:

\[
- \sum_{\alpha \in \{x, y, z\}} (J_{\rho})_{jl} + (L_{\rho})_{jl} + 2 \sum_{k \in V(i)} K_{ik} + 2 \sum_{k \in E(i)} K_{ik} + 2 \sum_{k \in K(i)} K_{ik} = 0.
\]

(30)

Setting the column vectors $\hat{\mathbf{V}} = (\hat{\mathbf{V}}_i)_{1 \leq i \leq N_T}$ and $\hat{\sigma} = (\hat{\sigma}_i)_{1 \leq i \leq N_T}$ whose sizes are $N_T \times \text{dof}$, the system (27)-(28) can be rewritten as:

\[
\begin{cases}
M_{\rho} \partial_t \hat{\mathbf{V}} = Q \hat{\sigma} - R \hat{\mathbf{V}}, \\
M_{\Lambda} \partial_t \hat{\sigma} = X \hat{\mathbf{V}} - Y \hat{\sigma},
\end{cases}
\]

(31)

where $M_{\rho}$, $M_{\Lambda}$, $R$ and $Y$ are $(N_T \times \text{dof}) \times (N_T \times \text{dof})$ block diagonal matrices such that for $1 \leq i \leq N_T$,

\[
(M_{\rho})_{ii} = M_{\rho_1}, \quad (M_{\Lambda})_{ii} = M_{\Lambda_i}, \quad R_{ii} = \sum_{k \in E(i)} \sigma_{ik}, \quad Y_{ii} = \sum_{k \in K(i)} \rho_{ik}.
\]

(32)

More precisely, $M_{\rho}$, $M_{\Lambda}$ and $R$ are symmetric definite positive matrices, and $Y$ is a symmetric semi-positive matrix. On the other hand, $Q$ and $X$ are $(N_T \times \text{dof}) \times (N_T \times \text{dof})$ block sparse matrices such that for $1 \leq i, j \leq N_T$,

\[
Q_{ij} = \begin{cases}
- \sum_{\alpha \in \{x, y, z\}} (J_{\rho})_{ij} + \sum_{k \in V(i)} K_{ik} + \sum_{k \in E(i)} K_{ik}, & \text{if } i = j, \\
G_{ij}, & \text{if } j \in V(i), \\
0, & \text{elsewhere},
\end{cases}
\]

(33)

and

\[
X_{ij} = \begin{cases}
- \sum_{\alpha \in \{x, y, z\}} (L_{\rho})_{ij} + \sum_{k \in V(i)} K_{ik} + \sum_{k \in E(i)} K_{ik} + 2 \sum_{k \in K(i)} K_{ik}, & \text{if } i = j, \\
- G_{ij}, & \text{if } j \in V(i), \\
0, & \text{elsewhere}.
\end{cases}
\]

(34)

Applying (30), the following essential property holds:

\[
Q + X^T = 0.
\]

(35)

### 2.4.2 High-order leap-frog scheme

Then, we apply a high-order leap-frog time-integration scheme. It is an explicit scheme which results, when combined to the centered fluxes defined at (14), in a non-diffusive scheme (see section 3.2). We note by $\Delta t$ the time-step and the superscripts refer to time stations. Thus, $\hat{\mathbf{V}}^n$ represents the velocity field at $t_n = n \Delta t$ and $\hat{\sigma}^{n+\frac{1}{2}}$ the stress tensor components at $t_{n+\frac{1}{2}} = (n + \frac{1}{2}) \Delta t$. Note that linear combinations (12) are still valid with superscripts. According to Young [62], the leap-frog integration scheme with an order $N$ (even) is written:

\[
\begin{pmatrix}
\hat{\mathbf{V}}(t_{n+1}) \\
\hat{\sigma}(t_{n+\frac{1}{2}})
\end{pmatrix} = \begin{pmatrix}
\hat{\mathbf{V}}(t_n) \\
\hat{\sigma}(t_{n+\frac{1}{2}})
\end{pmatrix} + 2 \sum_{m=1}^{N-1} \frac{1}{m!} \left( \frac{\Delta t}{2} \right)^m \partial_t^m \begin{pmatrix}
\hat{\mathbf{V}}(t_{n+\frac{1}{2}}) \\
\hat{\sigma}(t_{n+1})
\end{pmatrix}.
\]

(36)
As \( \hat{V} \) and \( \hat{\sigma} \) are not defined at \( t_{n+\frac{1}{2}} \) and \( t_{n+1} \) respectively, we introduce the linear combinations \( \hat{V}^{[n+\frac{1}{2}]} \) and \( \hat{\sigma}^{[n+1]} \): consider \( \alpha \in \mathbb{N}^* \), then we denote by

\[
\hat{V}^{[n+\frac{1}{2}]} = \sum_{i=1}^{\alpha} b_i \hat{V}^{n+1-i} \quad \text{and} \quad \hat{\sigma}^{[n+1]} = \sum_{i=1}^{\alpha} b_i \hat{\sigma}^{n+\frac{3}{2}-i},
\]

where the coefficients \( b_i \) satisfy the following linear Vandermonde system whose solution is unique:

\[
\sum_{i=1}^{\alpha} b_i = 1 \quad \text{and} \quad \sum_{i=1}^{\alpha} b_i \left( \frac{1}{2} - i \right)^k = 0, \quad \forall k \in \{1, ..., \alpha - 1\} \quad \text{when} \quad \alpha \geq 2.
\]

For example, we obtain for \( \alpha \in \{1, 2, 3, 4\} \):

\[
\begin{align*}
\hat{V}^{[n+\frac{1}{2}]} &= \hat{V}^n, \quad \hat{V}^{[n+\frac{1}{2}]} = \frac{3 \hat{V}^n - \hat{V}^{n-1}}{2}, \quad \hat{V}^{[n+\frac{1}{2}]} = \frac{15 \hat{V}^n - 10 \hat{V}^{n-1} + 3 \hat{V}^{n-2}}{8}, \\
\hat{V}^{[n+\frac{1}{2}]} &= 35 \hat{V}^n - 35 \hat{V}^{n-1} + 21 \hat{V}^{n-2} - 5 \hat{V}^{n-3}.
\end{align*}
\]

The interest of such a choice is that (see (135)-(136) for more details):

\[
\begin{align*}
\sum_{i=1}^{\alpha} b_i \hat{V}(t_{n+1-i}) &= \hat{V}(t_{n+\frac{1}{2}}) + \mathcal{O}(\Delta t^\alpha) \quad \text{and} \quad \sum_{i=1}^{\alpha} b_i \hat{\sigma}(t_{n+\frac{3}{2}-i}) = \hat{\sigma}(t_{n+1}) + \mathcal{O}(\Delta t^\alpha).
\end{align*}
\]

On the other hand, using (31), for all \( m \) (odd), we denote by \( \alpha_m, \beta_m, \gamma_m \) and \( \delta_m \) the matrices computed in the following way:

\[
\partial_t^m \left( \begin{array}{c} \hat{V} \\ \hat{\sigma} \end{array} \right) = \left( \begin{array}{cc} -M_r^{-1}R & M_r^{-1}Q \\ M_r^{-1}X & -M_r^{-1}Y \end{array} \right)^m \left( \begin{array}{c} \hat{V} \\ \hat{\sigma} \end{array} \right) = \left( \begin{array}{cc} \alpha_m & \beta_m \\ \gamma_m & \delta_m \end{array} \right) \left( \begin{array}{c} \hat{V} \\ \hat{\sigma} \end{array} \right).
\]

Starting from initial values \( \hat{V}_i^0 \) at \( t = 0 \) and \( \hat{\sigma}_i^\frac{1}{2} \) at \( t = \frac{\Delta t}{2} \), the final scheme is explicit and expresses:

\[
\begin{align*}
R_N &= -2M_r \sum_{m=1 \text{ (odd)}}^{N-1} \frac{1}{m!} \left( \frac{\Delta t}{2} \right)^m \alpha_m, \quad Q_N = 2M_r \sum_{m=1 \text{ (odd)}}^{N-1} \frac{1}{m!} \left( \frac{\Delta t}{2} \right)^m \beta_m, \\
X_N &= 2M_\Lambda \sum_{m=1 \text{ (odd)}}^{N-1} \frac{1}{m!} \left( \frac{\Delta t}{2} \right)^m \gamma_m, \quad Y_N = -2M_\Lambda \sum_{m=1 \text{ (odd)}}^{N-1} \frac{1}{m!} \left( \frac{\Delta t}{2} \right)^m \delta_m.
\end{align*}
\]

2.4.3 Properties

**Lemma 1** For all \( N \in \mathbb{N}^* \) (even), the matrices \( Q_N, X_N, R_N \) and \( Y_N \) satisfy the following properties:

\[
Q_N + X_N = 0, \quad R_N = R_N^t, \quad Y_N = Y_N^t.
\]

**Proof 1** We recursively prove that for all \( m \in \mathbb{N}^* \):

\[
\alpha_m^t M_r = M_r \alpha_m, \quad \delta_m^t M_\Lambda = M_\Lambda \delta_m, \quad M_r \beta_m + \gamma_m^t M_\Lambda = 0.
\]

For \( m = 1 \), as \( M_r, M_\Lambda, R, Y \) are symmetric, then (41) and (55) give:

\[
\alpha_1^t M_r = -R = M_r \alpha_1, \quad \delta_1^t M_\Lambda = -Y = M_\Lambda \delta_1, \quad M_r \beta_1 + \gamma_1^t M_\Lambda = Q + X^t = 0.
\]
After, we remark that
\[
\begin{pmatrix}
\alpha_{m+1} \\
\beta_{m+1} \\
\gamma_{m+1} \\
\delta_{m+1}
\end{pmatrix}
= \begin{pmatrix}
\alpha_1 & \beta_1 \\
\gamma_1 & \delta_1
\end{pmatrix}
\begin{pmatrix}
\alpha_m \\
\beta_m \\
\gamma_m \\
\delta_m
\end{pmatrix}
= \begin{pmatrix}
\alpha_1 & \beta_1 \\
\gamma_1 & \delta_1
\end{pmatrix}
\begin{pmatrix}
\alpha_m \\
\beta_m \\
\gamma_m \\
\delta_m
\end{pmatrix}.
\]

Therefore, according to (46), (47) and (48), we obtain
\[
\alpha_{m+1} M_\rho = (\alpha_1 \alpha_m + \beta_1 \gamma_m)^t M_\rho = \alpha^t M_\rho \alpha_1 M_\rho + \gamma^t M_\rho \beta_1 M_\rho
= M_\rho (\alpha_1 \alpha_m + \beta_1 \gamma_m) = M_\rho \alpha_m, \tag{49}
\]
\[
\delta_{m+1} M_\Lambda = (\gamma_1 \beta_m + \delta_1 \delta_m)^t M_\Lambda = \beta^t M_\Lambda \gamma_1 M_\Lambda + \delta^t M_\Lambda \delta_1 M_\Lambda
= M_\Lambda (\gamma_1 \beta_m + \delta_1 \delta_m) = M_\Lambda \delta_m, \tag{50}
\]
\[
M_\rho \beta_{m+1} + \gamma_{m+1} M_\Lambda = M_\rho (\alpha_1 \beta_m + \beta_m \delta_1) + (\gamma_1 \alpha_m + \delta_1 \gamma_m)^t M_\Lambda
= (M_\rho \alpha_1 \beta_m + \gamma_1 M_\Lambda) + (M_\rho \beta_1 \gamma_m + \gamma_1 M_\Lambda) \delta_1
= \alpha^t M_\rho (\beta_1 M_\Lambda) + (M_\rho \beta_m + \gamma_1 M_\Lambda) \delta_1
= 0. \tag{51}
\]

\[\square\]

Remark 3 In the particular case of free surface boundary conditions only (i.e. when \(R = 0\) and \(Y = 0\)), for \(m\) odd, we recursively show that:
\[
\alpha_m = 0, \quad \beta_m = M_\rho^{-1} Q \left( M_\Lambda^{-1} X M_\rho^{-1} Q \right)^{m-1} \gamma_m = M_\Lambda^{-1} X (M_\rho^{-1} Q M_\Lambda^{-1} X)^{m-1} \quad \text{and} \quad \delta_m = 0. \tag{52}
\]

Thus, with a change of variable, we find that \(R_N = 0, Y_N = 0\) and
\[
Q_N = \sum_{k=0}^{N-1} \frac{1}{(2k+1)!} \left( \frac{\Delta t}{2} \right)^{2k} Q \left( M_\Lambda^{-1} X M_\rho^{-1} Q \right)^k, \tag{53}
\]
\[
X_N = \sum_{k=0}^{N-1} \frac{1}{(2k+1)!} \left( \frac{\Delta t}{2} \right)^{2k} X \left( M_\rho^{-1} Q M_\Lambda^{-1} X \right)^k. \tag{54}
\]

Remark 4 As the matrix \(R_2 = R\) is symmetric positive definite and \(M_\rho \alpha_3\) is symmetric, then we always can find a condition under the time step \(\Delta t\) such that
\[
R_4 = R - \frac{\Delta t^2}{24} M_\rho \alpha_3 \tag{55}
\]
is a symmetric positive definite matrix.

Remark 5 As the matrix \(Y_2 = Y\) is symmetric positive semi-definite, \(M_\Lambda \delta_3\) is symmetric, and
\[
Y_4 = Y + \frac{\Delta t^2}{24} \left[ XM_\rho^{-1} RM_\rho^{-1} Q + YM_\Lambda^{-1} XM_\rho^{-1} Q + XM_\rho^{-1} Q M_\Lambda^{-1} Y + YM_\Lambda^{-1} Y M_\Lambda^{-1} Y \right] \tag{56}
\]
then by considering the eigenvector \(v \neq 0\) of \(Y\) such that \(Yv = 0\), we find \(v^t Y_4 v < 0\). Consequently, \(Y_4\) is only symmetric (but not positive semi-definite).

3 Energy preservation and stability of the scheme

The discrete scheme (42) being explicit, it is conditionally stable. The aim of this section is to establish a sufficient condition on the time step \(\Delta t\) for the \(L^2\)-stability of the Discontinuous Galerkin scheme, taking into account both free surface (15) and absorbing boundary conditions (23)-(24). Following [30, 12, 13], we shall define a discrete energy and prove that it is a quadratic positive definite form playing the role of a Lyapunov function of numerical unknowns.
3.1 A discrete elastodynamic energy

In the continuum, the total energy of the system is given by:

\[ E(t) = \frac{1}{2} \int_{\Omega} \rho \| \vec{V}(t) \|^2 \, dv + \frac{1}{2} \int_{\Omega} \vec{\sigma}(t) : \epsilon(t) \, dv, \]

where \( \vec{\sigma} \) is the symmetric deformation tensor whose components are \( \epsilon_{ij} = \frac{1}{2} \left( \frac{\partial U_i}{\partial x_j} + \frac{\partial U_j}{\partial x_i} \right) \). As for the stress, we introduce the deformation vector \( \vec{\epsilon} = (\epsilon_{xx}, \epsilon_{yy}, \epsilon_{zz}, 2\epsilon_{xy}, 2\epsilon_{xz}, 2\epsilon_{yz})^t \). For an elastic medium, the generalized Hooke’s law links deformations and stresses through the linear relationship \( \vec{\epsilon} = C \vec{\sigma} \), where the matrix \( C \) expresses:

\[
C = \frac{1}{E} \begin{pmatrix}
1 & -\nu & -\nu & 0 & 0 & 0 \\
-\nu & 1 & -\nu & 0 & 0 & 0 \\
-\nu & -\nu & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & 2(1+\nu) & 0 & 0 \\
0 & 0 & 0 & 0 & 2(1+\nu) & 0 \\
0 & 0 & 0 & 0 & 0 & 2(1+\nu)
\end{pmatrix},
\]

where \( E \) and \( \nu \) being respectively the Young’s modulus and the Poisson’s ratio

\[ E = \frac{\mu (3\lambda + 2\mu)}{\lambda + \mu} \quad \text{and} \quad \nu = \frac{\lambda}{2(\lambda + \mu)}. \]

Thanks to this relationship, the mechanical energy becomes

\[ E_m = \frac{1}{2} \int_{\Omega} (\vec{\sigma})^t \epsilon \, dv = \frac{1}{2} \int_{\Omega} (\vec{\sigma})^t C \vec{\sigma} \, dv, \]

and introducing the change of variables (4), we obtain

\[ E_m = \frac{1}{2} \int_{\Omega} (T^{-1} \vec{\sigma})^t C (T^{-1} \vec{\sigma}) \, dv = \frac{1}{2} \int_{\Omega} (\vec{\theta})^t \left[(T^{-1})^t C T^{-1}\right] \vec{\theta} \, dv. \]

We can easily check that \( (T^{-1})^t C T^{-1} = \Lambda \) so that, the total energy writes

\[ E(t) = \frac{1}{2} \int_{\Omega} \rho \| \vec{V}(t) \|^2 \, dv + \frac{1}{2} \int_{\Omega} (\vec{\theta}(t))^t \Lambda(\lambda, \mu) \vec{\theta}(t) \, dv. \]

**Definition 1** The discrete energy \( E^n \) in three dimensions of space at time \( n\Delta t \) is defined by:

\[
E^n = \frac{1}{2} \left[ \left( \hat{\vec{V}}^n \right)^t M_\rho \hat{\vec{V}}^n + \left( \hat{\vec{\sigma}}^{n+\frac{1}{2}} \right)^t \Lambda \hat{\vec{\sigma}}^{n-\frac{1}{2}} \right].
\]  
(57)

3.2 Energy preservation

Now, we study the evolution of the discrete elastodynamic energy through one time-step. More precisely, we aim to establish that the combination between the centered approximation of the fluxes and the high-order leap-frog time-integration leads to a non-dissipative scheme for unbounded domains or domains with free surfaces.

**Lemma 2** For all \( N \geq 2 \) (even), the variation of the total discrete energy during one time-step \( \Delta t \) is

\[
E^{n+1} - E^n = -\Delta t \left( \frac{\hat{\vec{V}}^{n+1} + \hat{\vec{V}}^n}{2} \right)^t R_N \hat{\vec{V}}^{[n+\frac{1}{2}]} - \Delta t \left( \hat{\vec{\sigma}}^{n+\frac{1}{2}} \right)^t Y_N \frac{\hat{\vec{\sigma}}^{n+1} + \hat{\vec{\sigma}}^n}{2}.
\]  
(58)
Proof 2 We calculate the variation of the discrete energy during one time-step $\Delta t$:

$$E^{n+1} - E^n = \frac{1}{2}M_{\rho} \left( \hat{\nabla}^{n+1} - \hat{\nabla}^n \right) + \frac{1}{2} M_{\lambda} \left( \hat{\sigma}^{n+\frac{1}{2}} - \hat{\sigma}^{n-\frac{1}{2}} \right).$$

(59)

Using (42), we find

$$E^{n+1} - E^n = \Delta t \left( \frac{\hat{\nabla}^{n+1} + \hat{\nabla}^n}{2} \right)^t \left( (Q_N + X_N^t) \hat{\sigma}^{n+\frac{1}{2}} - R_N \hat{\nabla}^{n+\frac{1}{2}} \right) - \Delta t \left( \hat{\sigma}^{n+\frac{1}{2}} \right)^t \left( \frac{\hat{\sigma}^{n+1} + \hat{\sigma}^n}{2} \right).$$

(60)

We conclude with (45). □

Corollary 1 Using the scheme (42), for an infinite domain or a domain including free surface boundaries (15) only, the discrete elastodynamic energy is preserved through one time-step:

$$\forall n \in \mathbb{N}, \ E^{n+1} = E^n.$$  

Proof 3 In an infinite domain or a domain with free surface boundary conditions, $\mathcal{E}(i)$ is an empty set, so that the matrices $R = 0$ and $Y = 0$. Therefore, following remark 3 and (58), the discrete elastodynamic energy is preserved. □

3.3 A corrected discrete elastodynamic energy for absorbing boundaries

Let be $\alpha \in \mathbb{N}^*$ and $b_i$ as in (37), $\forall i \in \{1, ..., \alpha\}$. In what follows, we define the signs of $b_i$ and $b_i + b_{i+1}$ respectively by:

$$\forall i \in \{1, ..., \alpha\}, \ \gamma_i = \frac{b_i}{|b_i|},$$

(61)

$$\forall i \in \{1, ..., \alpha - 1\}, \ s_i = \frac{b_i + b_{i+1}}{|b_i + b_{i+1}|}.$$  

(62)

When absorbing boundary conditions are applied at some faces of the domain (i.e. when $\mathcal{E}(i)$ is not empty), the discrete variation of the energy through one time-step $\Delta t$ is not necessarily negative according to Lemma 2, because of the time-asymmetry on both $\hat{V}$ and $\hat{\sigma}$ in (58). In order to overcome this difficulty, we introduce correction terms in the definition of the discrete energy (57) and prove that this corrected discrete elastodynamic energy is not increased through each time step.

Definition 2 For $\alpha = 1$, we introduce a corrected discrete elastodynamic energy $\mathcal{E}^n$ at time $n\Delta t$ by:

$$\mathcal{E}^n = E^n - \frac{\Delta t}{4} \left[ \left( \hat{\nabla}^n \right)^t R_N \hat{\nabla}^n - \left( \hat{\sigma}^{n-\frac{1}{2}} \right)^t Y_N \hat{\sigma}^{n-\frac{1}{2}} \right].$$

(63)

Definition 3 For $\alpha \in \{2, 3, 4\}$, we introduce a corrected discrete elastodynamic energy $\mathcal{E}^n$ at time $n\Delta t$ by:

$$\mathcal{E}^n = E^n - \frac{\Delta t}{4} \sum_{i=1}^{\alpha} |b_i| \left[ \left( \hat{\nabla}^n \right)^t R_N \hat{\nabla}^n - \left( \hat{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \hat{\sigma}^{n+\frac{1}{2} - i} \right]$$

$$+ \frac{\Delta t}{4} \sum_{i=2}^{\alpha} |b_i| \left( \hat{\nabla}^n - \gamma_i \hat{\nabla}^{n+1-i} \right)^t R_N \left( \hat{\nabla}^n - \gamma_i \hat{\nabla}^{n+1-i} \right).$$

(64)

These correction terms only concern the absorbing faces and have no particular physical meaning. However, they match the loss of energy (which appears as a decreasing discrete energy) through these faces since we only consider outgoing waves. Using (58), we easily prove the following lemma when $\alpha = 1$.

Lemma 3 For $\alpha = 1$, we have for all $n \in \mathbb{N}$:

$$\mathcal{E}^{n+1} - \mathcal{E}^n = -\frac{\Delta t}{4} \left[ \left( \hat{\nabla}^{n+1} + \hat{\nabla}^n \right)^t R_N \left( \hat{\nabla}^{n+1} + \hat{\nabla}^n \right) + \left( \hat{\sigma}^{n+\frac{1}{2}} + \hat{\sigma}^{n-\frac{1}{2}} \right)^t Y_N \left( \hat{\sigma}^{n+\frac{1}{2}} + \hat{\sigma}^{n-\frac{1}{2}} \right) \right].$$

(65)
Then, we can establish the following result when \( \alpha \in \{2, 3, 4\} \), but it is not valid for \( \alpha \geq 5 \).

**Lemma 4** For \( \alpha \in \{2, 3, 4\} \), we have for all \( n \in \mathbb{N} \):

\[
E^{n+1} - E^n = \frac{-\Delta t}{4} \left[ \sum_{i=1}^{\alpha-1} [b_i + b_{i+1}] \left( \hat{\nu}^{n+1} + s_i \hat{\nu}^{n+1-i} \right)^t R_N \left( \hat{\nu}^{n+1} + s_i \hat{\nu}^{n+1-i} \right) \right. \\
+ |b_1| \left( \hat{\nu}^{n+1} + \gamma_1 \hat{\nu}^{n+1-i} \right)^t R_N \left( \hat{\nu}^{n+1} + \gamma_1 \hat{\nu}^{n+1-i} \right) \\
+ \left. \sum_{i=1}^{\alpha-1} |b_i + b_{i+1}| \left( \tilde{\sigma}^{n+\frac{1}{2}} + s_i \tilde{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \left( \tilde{\sigma}^{n+\frac{1}{2}} + s_i \tilde{\sigma}^{n+\frac{1}{2} - i} \right) \right]\). \tag{66}

**Proof 4** We denote by \( E_V^n = \frac{1}{2} \left( \hat{\nu}^n \right)^t M \hat{\nu}^n \) and \( E_\sigma^n = \frac{1}{2} \left( \tilde{\sigma}^{n+\frac{1}{2}} \right)^t M \tilde{\sigma}^{n+\frac{1}{2}} \). Thus, according to (57)-(58) and the definitions (37) of \( \hat{\nu}^{[n+\frac{1}{2}]} \) and \( \tilde{\sigma}^{[n+1]} \), we obtain

\[
E_{V}^{n+1} - E_{V}^n = -\frac{\Delta t}{2} \left[ \sum_{i=1}^{\alpha} b_i \left( \hat{\nu}^{n+1} \right)^t R_N \hat{\nu}^{n+1-i} + b_1 \left( \hat{\nu}^n \right)^t R_N \hat{\nu}^n + \sum_{i=2}^{\alpha} b_i \left( \hat{\nu}^n \right)^t R_N \hat{\nu}^{n+1-i} \right] \tag{67}
\]

\[
E_\sigma^{n+1} - E_\sigma^n = -\frac{\Delta t}{2} \left( \tilde{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \left( b_1 \tilde{\sigma}^{n+\frac{1}{2}} + \sum_{i=1}^{\alpha-1} (b_i + b_{i+1}) \tilde{\sigma}^{n+\frac{1}{2} - i} + b_\alpha \tilde{\sigma}^{n+\frac{1}{2} - \alpha} \right) \tag{68}
\]

Then, we set \( \mathcal{E} = \mathcal{E}_V + \mathcal{E}_\sigma \), with

\[
\mathcal{E}_V^n = E_V^n - \frac{\Delta t}{4} \left[ |b_1| \left( \hat{\nu}^n \right)^t R_N \hat{\nu}^n - \sum_{i=2}^{\alpha} |b_i| \left( \hat{\nu}^{n+1-i} \right)^t R_N \hat{\nu}^{n+1-i} + 2 \sum_{i=2}^{\alpha} b_i \left( \hat{\nu}^n \right)^t R_N \hat{\nu}^{n+1-i} \right] \tag{69}
\]

\[
\mathcal{E}_\sigma^n = E_\sigma^n + \frac{\Delta t}{4} \sum_{i=1}^{\alpha} |b_i| \left( \tilde{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2} - i} \tag{70}
\]

so that, according to (67)-(69) and (68)-(70), we have

\[
\mathcal{E}_{V}^{n+1} - \mathcal{E}_{V}^n = -\frac{\Delta t}{4} \left[ |b_1| \left( \hat{\nu}^{n+1} \right)^t R_N \hat{\nu}^{n+1} + (2b_1 - |b_1| - |b_2|) \left( \hat{\nu}^n \right)^t R_N \hat{\nu}^n \\
+ \sum_{i=2}^{\alpha-1} (|b_i| - |b_{i+1}|) \left( \hat{\nu}^{n+1-i} \right)^t R_N \hat{\nu}^{n+1-i} + |b_\alpha| \left( \hat{\nu}^{n+1-\alpha} \right)^t R_N \hat{\nu}^{n+1-\alpha} \\
+ 2 \sum_{i=1}^{\alpha-1} (b_i + b_{i+1}) \left( \hat{\nu}^{n+1} \right)^t R_N \hat{\nu}^{n+1-i} + 2b_\alpha \left( \hat{\nu}^{n+1} \right)^t R_N \hat{\nu}^{n+1-\alpha} \right] \tag{71}
\]

\[
\mathcal{E}_\sigma^{n+1} - \mathcal{E}_\sigma^n = -\frac{\Delta t}{4} \left[ -|b_1| \left( \tilde{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2}} + \sum_{i=1}^{\alpha-1} (|b_i| - |b_{i+1}|) \left( \tilde{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2} - i} \\
+ |b_\alpha| \left( \tilde{\sigma}^{n+\frac{1}{2} - \alpha} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2} - \alpha} \right] + E_{V}^{n+1} - E_{V}^n \tag{72}
\]

At last, we use the following identities in (68): for \( i \in \{1, ..., \alpha - 1\} \),

\[
2(b_i + b_{i+1}) \left( \tilde{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2} - i} = \left[ b_i + b_{i+1} \left( \tilde{\sigma}^{n+\frac{1}{2}} + s_i \tilde{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \left( \tilde{\sigma}^{n+\frac{1}{2}} + s_i \tilde{\sigma}^{n+\frac{1}{2} - i} \right) \\
- \left[ b_i + b_{i+1} \left( \tilde{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2}} + \left( \tilde{\sigma}^{n+\frac{1}{2} - i} \right)^t Y_N \tilde{\sigma}^{n+\frac{1}{2} - i} \right] \right] \tag{73}
\]
Proof 5 We use lemmas 3 and 4. If \( \| E \| \) and similar identities for the two last terms of (71). We finally obtain the variation of the total corrected discrete elastodynamic energy:

\[
E^{n+1} - E^n = \frac{-\Delta t}{4} \left[ 2(b_1 - |b_1|) \left( \hat{\gamma}^n \right)^t R_N \hat{\gamma}^n + \left( |b_1| - \sum_{i=1}^{\alpha - 1} |b_i + b_{i+1}| - |b_{\alpha}| \right) \left( \hat{\gamma}^{n+1} \right)^t R_N \hat{\gamma}^{n+1} \right]
\]

\[
+ \left( 2b_1 - |b_1| - \sum_{i=1}^{\alpha - 1} |b_i + b_{i+1}| - |b_{\alpha}| \right) \left( \hat{\sigma}^{n+1-i} \right)^t Y_N \hat{\sigma}^{n+1-i}
\]

\[
+ \sum_{i=1}^{\alpha - 1} |b_i + b_{i+1}| \left( \hat{\gamma}^{n+1} + s_i \hat{\gamma}^{n+1-i} \right)^t R_N \left( \hat{\gamma}^{n+1} + s_i \hat{\gamma}^{n+1-i} \right) \]

\[
+ |b_{\alpha}| \left( \hat{\gamma}^{n+1} + \gamma \hat{\gamma}^{n+1-\alpha} \right)^t R_N \left( \hat{\gamma}^{n+1} + \gamma \hat{\gamma}^{n+1-\alpha} \right)
\]

\[
+ \sum_{i=1}^{\alpha - 1} |b_i + b_{i+1}| \left( \hat{\sigma}^{n+1+i} + s_i \hat{\sigma}^{n+1-i} \right)^t Y_N \left( \hat{\sigma}^{n+1+i} + s_i \hat{\sigma}^{n+1-i} \right)
\]

\[
+ |b_{\alpha}| \left( \hat{\sigma}^{n+1+i} + \gamma \hat{\sigma}^{n+1+i-\alpha} \right)^t Y_N \left( \hat{\sigma}^{n+1+i} + \gamma \hat{\sigma}^{n+1+i-\alpha} \right)
\]

(74)

When \( \alpha \in \{2, 3, 4\} \), we can verify according to (39) that \( b_1 \geq 0 \), then \( b_1 - \sum_{i=1}^{\alpha - 1} |b_i + b_{i+1}| - |b_{\alpha}| = 0 \) and \( \forall i \in \{1, ..., \alpha - 1\}, |b_i| - |b_i + b_{i+1}| - |b_{i+1}| = 0 \). On the contrary, these two last equalities are not satisfied for \( \alpha \in \{5, 6\} \).

\[ \Box \]

Corollary 2 For \( \alpha \in \{1, 2, 3, 4\} \) and \( N = 2 \), the sequel \( (E^n)_{n \in \mathbb{N}} \) is not increasing.

Proof 5 We use lemmas 3 and 4. If \( N = 2 \), then \( E^{n+1} - E^n \) is a negative quantity, because the matrices \( R \) and \( Y \) are positive definite symmetric and semi-positive definite symmetric respectively. On the contrary, if \( N \geq 4 \), we are unable to conclude according to remarks 4 and 5.

\[ \Box \]

3.4 Stability of the scheme

We aim at proving that the discrete corrected elastodynamic energy \( E^n \) is a positive definite quadratic form of the unknowns \( \hat{\gamma}^n \) and \( \hat{\sigma}^{n-1/2} \) under some stability condition on the time step \( \Delta t \). For that, we need to introduce another formulation of \( E^n \), independently of the unknowns \( \hat{\sigma}^{n+1/2} \).

Theorem 1 Considering the scheme (42), the corrected total discrete elastodynamic energy \( E^n \) is non-increasing through iterations and is a positive definite quadratic form of numerical unknowns under the condition:

1. for domains with free surface boundary conditions only, the discrete scheme (42) is \( L^2 \)-stable for all \( N \) (even) if

\[ \Delta t \leq \frac{2}{\| M_A^{-1/2} X N M_{\rho}^{-1/2} \|}, \]

(75)

2. elsewhere, the discrete scheme (42) is \( L^2 \)-stable for \( N = 2 \) and \( \alpha \in \{1, 2, 3, 4\} \) if

\[ \Delta t \leq \frac{2}{\| M_A^{-1/2} X M_{\rho}^{-1/2} \| + \left( \sum_{i=1}^{\alpha} |b_i| \right) \max \left( \| M_{\rho}^{-1/2} R M_{\rho}^{-1/2} \|, \| M_A^{-1/2} Y M_A^{-1/2} \| \right)}, \]

(76)

where \( \| \cdot \| \) denotes a matricial norm.
Proof 6 With (42) in (57), it holds
\[ 2E_n = (\hat{\nu}_n)^t M_\rho \hat{\nu}_n + (\hat{\sigma}^{n-\frac{1}{2}})^t \left[ M_\Lambda \hat{\sigma}^{n-\frac{1}{2}} + \Delta t X_N \hat{\nu}_n - \Delta t Y_N \hat{\sigma}^{[n]} \right]. \] (77)

As $M_\rho$ and $M_\Lambda$ are symmetric definite positive matrices, we can bound down the discrete energy in the following way by using (57):
\[
2E_n \geq \left\| M_\rho^{1/2} \hat{\nu}_n \right\|^2 + \left\| M_\Lambda^{1/2} \hat{\sigma}^{n-\frac{1}{2}} \right\|^2 - \Delta t \left\| M_\Lambda^{1/2} \hat{\sigma}^{n-\frac{1}{2}} \right\| \left\| M_\rho^{-1/2} X N M_\rho^{-1/2} \right\| \left\| M_\rho^{1/2} \hat{\nu}_n \right\|
- \Delta t \sum_{i=1}^\alpha b_i \left( \hat{\sigma}^{n-\frac{1}{2}} \right)^t Y_N \hat{\sigma}^{n+\frac{1}{2}},
\] (78)

which gives the following estimate
\[
2E_n \geq \left\| M_\rho^{1/2} \hat{\nu}_n \right\|^2 \left( 1 - \frac{\Delta t}{2} \left\| M_\Lambda^{-1/2} X N M_\rho^{-1/2} \right\| \right) + \left\| M_\Lambda^{1/2} \hat{\sigma}^{n-\frac{1}{2}} \right\|^2 \left( 1 - \frac{\Delta t}{2} \left\| M_\Lambda^{-1/2} X N M_\rho^{-1/2} \right\| \right)
+ \frac{\Delta t}{2} \sum_{i=1}^\alpha |b_i| \left( \hat{\sigma}^{n-\frac{1}{2}} - \gamma_i \hat{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \left( \hat{\sigma}^{n-\frac{1}{2}} - \gamma_i \hat{\sigma}^{n+\frac{1}{2}} \right)
- \frac{\Delta t}{2} \sum_{i=1}^\alpha |b_i| \left( \hat{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \hat{\sigma}^{n+\frac{1}{2}}.
\] (80)

Then, we can bound down the discrete corrected energy from (64)
\[
\epsilon_n \geq \frac{1}{2} \left\| M_\rho^{1/2} \hat{\nu}_n \right\|^2 \left( 1 - \frac{\Delta t}{2} \left\| M_\Lambda^{-1/2} X N M_\rho^{-1/2} \right\| \right) - \frac{\Delta t}{2} \left( \sum_{i=1}^\alpha |b_i| \right) \left\| M_\rho^{-1/2} R_N M_\rho^{-1/2} \right\|
+ \frac{1}{2} \left\| M_\Lambda^{1/2} \hat{\sigma}^{n-\frac{1}{2}} \right\|^2 \left( 1 - \frac{\Delta t}{2} \left\| M_\Lambda^{-1/2} X N M_\rho^{-1/2} \right\| \right) - \frac{\Delta t}{2} \left( \sum_{i=1}^\alpha |b_i| \right) \left\| M_\Lambda^{-1/2} Y_N M_\Lambda^{-1/2} \right\|
+ \frac{\Delta t}{4} \sum_{i=1}^\alpha |b_i| \left( \hat{\nu}_n - \gamma_i \hat{\nu}_{n+1} \right)^t \left( \hat{\nu}_n - \gamma_i \hat{\nu}_{n+1} \right)
+ \frac{\Delta t}{4} \sum_{i=1}^\alpha |b_i| \left( \hat{\sigma}^{n-\frac{1}{2}} - \gamma_i \hat{\sigma}^{n+\frac{1}{2}} \right)^t Y_N \left( \hat{\sigma}^{n-\frac{1}{2}} - \gamma_i \hat{\sigma}^{n+\frac{1}{2}} \right)
\] (81)

Finally, we conclude with remarks 3.4-5 and corollaries 1 and 2. □

4 Convergence analysis

The objective of this section is to prove the convergence of the totally discretized scheme (42) following the method in [30, 56]. We consider a family of unstructured tetrahedral meshes $\mathcal{T}_h$, where $h = \max_{i \in [1,N_T]} h_i$, which forms a partition of the domain $\Omega$, i.e. $\Omega = \cup_{\mathcal{T}_i \in \mathcal{T}_h} \mathcal{T}_i$. We assume that the unstructured meshes $\mathcal{T}_h$ are uniformly shape regular in the sense that there exists a constant $\xi > 0$ such that
\[ \forall h > 0, \quad \forall \mathcal{T}_i \in \mathcal{T}_h : \quad \frac{h_i}{d_i} \leq \xi, \] (82)

where $d_i$ is the diameter of the biggest ball included in $\mathcal{T}_i$. Moreover, we assume that there exists $\eta > 0$ (independent of $h$) such that
\[ \forall h > 0, \quad \forall \mathcal{T}_i \in \mathcal{T}_h, \quad \forall k \in \mathcal{V}(i) : \quad \frac{h_k}{h_i} \leq \eta. \] (83)

In what follows, we denote by
\[
X_1^h = \{ \varphi_h \in L^2(\Omega)^3 : \forall i, \varphi_{h_{|T_i}} \in P_k(T_i)^3 \},
\] (84)
\[
X_2^h = \{ \psi_h \in L^2(\Omega)^6 : \forall i, \psi_{h_{|T_i}} \in P_k(T_i)^6 \},
\] (85)
the sets of discontinuous piecewise polynomial fields of degree at most \( k \) on each \( T_i \). We next introduce the broken Sobolev spaces

\[
PH^s(\Omega) = \left\{ v : \forall i, \ v|_{T_i} \in H^s(T_i) \right\}
\]

equipped with the norm \( \|v\|_{PH^s(\Omega)} = \left( \sum_{i=1}^{N_T} \|v|_{T_i}\|_{H^s(T_i)} \right)^{1/2} \), where \( \| \cdot \|_{H^s(T_i)} \) denotes the standard \( H^s \)-norm on \( T_i \).

### 4.1 Definition and properties of the semi-discretized scheme

First, we are interested in the study of consistency and stability of the spatially semi-discretized problem below. In what follows, we set \( \tilde{\varphi}_i \) (resp. \( \tilde{\psi}_i \)) the restriction of \( \varphi \in H^s(T_h)^3 \) (resp. \( \psi \in H^s(T_h)^6 \)) to the element \( T_i \) with \( s > \frac{1}{2} \). The semi-discrete solution \( (\tilde{V}_h, \tilde{\sigma}_h) \) defined in \( C \left( [0, T], H^1_k \times H^2_k \right) \) is the solution of the weak formulation (25)-(26), for any test field \( (\varphi_h, \psi_h) \in X_1^h \times X_2^h \), for \( 0 \leq t \leq T \) and all \( i \), with the initial values

\[
\tilde{V}_h(0) = P^1_h(\tilde{V}_0) \quad \text{and} \quad \tilde{\sigma}_h(0) = P^2_h(\tilde{\sigma}_0),
\]

where \( P^1_h : L^2(\Omega)^3 \rightarrow X_1^h \) and \( P^2_h : L^2(\Omega)^6 \rightarrow X_2^h \) denote respectively the orthogonal projection onto \( X_1^h \) and \( X_2^h \).

**Definition 4** For given vector fields \( \tilde{W} = (\tilde{V}, \tilde{\sigma}) \) and \( \tilde{W}' = (\tilde{V}', \tilde{\sigma}') \), we introduce the following bilinear forms which are well defined on \( X_1^h \times X_2^h \):

\[
m(\tilde{W}, \tilde{W}') = \sum_{i=1}^{N_T} \int_{T_i} \left[ p_i (\tilde{V}'_i)^t \tilde{V}_i + (\tilde{\sigma}'_i)^t \Lambda \tilde{\sigma}_i \right] dv
\]

\[
a(\tilde{W}, \tilde{W}') = \sum_{i=1}^{N_T} \int_{T_i} \sum_{\alpha \in \{x,y,z\}} \left[ (\partial_{\alpha} \tilde{V}'_i)^t M_{\alpha} \tilde{V}_i + (\partial_{\alpha} \tilde{\sigma}'_i)^t M_{\alpha} \tilde{\sigma}_i \right] dv
\]

\[
b(\tilde{W}, \tilde{W}') = -\sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \left[ (\tilde{V}'_i)^t P_{ik} \tilde{\sigma}_i + (\tilde{\sigma}'_i)^t P_{ik} \tilde{V}_i \right] ds + \frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \left[ \left( \tilde{V}'_i \right)^t P_{ik} \tilde{V}_i - \frac{1}{2} \left( \tilde{\sigma}'_i \right)^t P_{ik} \tilde{\sigma}_i \right] ds
\]

Thanks to (25) and (26), we remark the following result:

**Proposition 1** The semi-discrete solution \( \tilde{W}_h = (\tilde{V}_h, \tilde{\sigma}_h) \) of (25)-(26)-(87) satisfies

\[
m(\partial_t \tilde{W}_h, \tilde{W}'_h) + a(\tilde{W}_h, \tilde{W}'_h) + b(\tilde{W}_h, \tilde{W}'_h) = 0, \quad \forall \tilde{W}'_h \in X_1^h \times X_2^h.
\]

**Proposition 2** The exact solution \( \tilde{W} = (\tilde{V}, \tilde{\sigma}) \) of (6)-(7)-(8)-(9) satisfies the following property

\[
m(\partial_t \tilde{W}_h, \tilde{W}'_h) + a(\tilde{W}_h, \tilde{W}'_h) + b(\tilde{W}_h, \tilde{W}'_h) = 0, \quad \forall \tilde{W}'_h \in X_1^h \times X_2^h.
\]

**Proof 7** According to (10a)-(10b), we have to verify that the quantity \( b(\tilde{W}, \tilde{W}'_h) \) is equal to

\[
- \sum_{i=1}^{N_T} \sum_{\alpha \in \{x,y,z\}} \int_{\partial T_i} \left[ (\tilde{V}'_i)^t M_{\alpha} \tilde{\sigma}_n \alpha + (\tilde{\sigma}'_i)^t M_{\alpha} \tilde{V}_n \alpha \right] ds.
\]

The free surface boundary condition (7) results in:

\[
\sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} (\tilde{V}'_i)^t P_{ik} \tilde{\sigma} ds = 0.
\]

On the other hand, thanks to the absorbing boundary conditions (8), we also have:

\[
\frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \left[ (\tilde{V}'_i)^t (P_{ik} \tilde{\sigma} - A_{ik} \tilde{V}) + (\tilde{\sigma}'_i)^t (P_{ik} \tilde{V} - B_{ik} \tilde{\sigma}) \right] ds = \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_{ik}} \left[ (\tilde{V}'_i)^t P_{ik} \tilde{\sigma} + (\tilde{\sigma}'_i)^t P_{ik} \tilde{V} \right] ds.
\]
which closes the proof.

Lemma 5 For any \( \tilde{W}_h = (\tilde{\varphi}_h, \tilde{\psi}_h) \in X_h^1 \times X_h^2 \), the following quantity is positive:

\[
a(W_h, W_h) + b(W_h, W_h) = \frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \left[ \varphi_i^t A_{ik} \varphi_i + \psi_i^t B_{ik} \psi_i \right] ds \geq 0. \tag{94}
\]

Proof 8 Let be \( \tilde{W}_h = (\tilde{\varphi}_h, \tilde{\psi}_h) \in X_h^1 \times X_h^2 \) a given vector field. Then, we get

\[
a(W_h, W_h) + b(W_h, W_h) = \sum_{i=1}^{N_T} \sum_{\alpha \in \{x,y,z\}} \int_{T_i} \left[ (\partial_\alpha \tilde{\varphi}_i)^t M_{\alpha} \tilde{\psi}_i + (\partial_\alpha \tilde{\psi}_i)^t M_{\alpha}^t \tilde{\varphi}_i \right] dv
- \frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \left[ \varphi_i^t P_{ik} \varphi_i + \psi_i^t P_{ik}^t \psi_i \right] ds
- \frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \left[ \varphi_i^t P_{ik} \varphi_k + \psi_i^t P_{ik}^t \varphi_k \right] ds
+ \frac{1}{2} \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \left[ \varphi_i^t A_{ik} \varphi_i + \psi_i^t B_{ik} \psi_i \right] ds. \tag{95}
\]

The first term (between brackets) vanishes by the Green Formula:

\[
\int_{T_i} \left[ (\partial_\alpha \tilde{\varphi}_i)^t M_{\alpha} \tilde{\psi}_i + (\partial_\alpha \tilde{\psi}_i)^t M_{\alpha}^t \tilde{\varphi}_i \right] dv = \int_{\partial T_i} \tilde{\varphi}_i^t M_{\alpha} n_\alpha \tilde{\psi}_i ds,
\]

combined to the equality \( \tilde{\psi}_i^t P_{ik}^t \tilde{\varphi}_i = \varphi_i^t P_{ik} \varphi_i \). Moreover, since \( P_{ik} = \sum_{\alpha \in \{x,y,z\}} M_{\alpha} n_{\alpha ik} \) and thanks to the orientation of \( \bar{u}_{ik} \) from \( T_i \) to \( T_k \), we remark that:

\[
\sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \psi_i^t P_{ik} \varphi_k ds = - \sum_{i=1}^{N_T} \sum_{k \in \mathcal{E}(i)} \int_{S_k} \varphi_i^t P_{ik} \psi_k ds,
\]

which implies that the second term is also equal to zero. Therefore, it remains only the third term depending on \( \bar{A}_{ik} \) and \( \bar{B}_{ik} \), which are respectively a real symmetric definite positive matrix and a real symmetric semi-definite positive matrix. \( \square \)

Consequently, if we define the semi-discrete energy by

\[
\mathcal{E}_h(t) = \frac{1}{2} m(W_h(t), W_h(t)), \tag{96}
\]

where \( W_h(t) = (\bar{V}_h(t), \bar{\sigma}_h(t)) \) is the semi-discrete solution of (25)-(26)-(87), we deduce from Prop. 1 and Lemma 5 that

\[
\partial_t \mathcal{E}_h(t) \leq 0, \quad \forall t \in [0,T], \tag{97}
\]

which implies that \( \mathcal{E}_h(t) \) is decreasing on \([0,T]\) and we have

\[
\mathcal{E}_h(t) \leq \mathcal{E}_h(0) \leq E(0), \quad \forall t \in [0,T]. \tag{98}
\]
4.2 Convergence of the semi-discretized problem

We shall at first give a convergence result for the spatially semi-discretized scheme (25)-(26)-(87). We recall the two following lemmata, well-known in the framework of the finite elements (see [17]):

**Lemma 6** Let $\mathcal{T}_i \in \mathcal{T}_h$, we assume that $u$ belongs to the space $H^{s+1}(\mathcal{T}_i)$ for $s \geq 0$. Let $\Pi$ be a linear continuous operator from $H^{s+1}(\mathcal{T}_i)$ onto $P_k(\mathcal{T}_i)$ such that $\Pi(u) = u$ for all $u \in P_k(\mathcal{T}_i)$. Then we have

$$
\|u - \Pi(u)\|_{L^2(\mathcal{T}_i)} \leq C_1 h_i^{\min\{s,k\}+1} \|u\|_{H^{s+1}(\mathcal{T}_i)},
$$

$$
\|u - \Pi(u)\|_{L^2(\partial\mathcal{T}_i)} \leq C_1 h_i^{\min\{s,k\}+1/2} \|u\|_{H^{s+1}(\mathcal{T}_i)},
$$

where $C_1$ is a positive constant only depending on $k$, $s$ and the regularity parameter $\xi$ of the mesh.

**Lemma 7** For all $p \in P_k(\mathcal{T}_i)$, we have

$$
\|p\|_{L^2(\partial\mathcal{T}_i)} \leq C_2 h_i^{-1/2} \|p\|_{L^2(\mathcal{T}_i)},
$$

where $C_2$ is a positive constant only depending on $k$ and the regularity parameter $\xi$ of the mesh.

Thus, we can state the following Lemma:

**Lemma 8** Let be $\tilde{W} = (\tilde{V}, \tilde{\sigma})$ the exact solution of (6)-(7)-(8)-(9) supposed to belong to $C^0([0,T], PH^{s+1}(\Omega)^9)$ and $(\tilde{V}_h, \tilde{\sigma}_h) \in C^1([0,T], X^1_h \times X^2_h)$ be the solution of the semi-discrete problem (25)-(26)-(87). For any $t \in [0,T]$, we have

$$
b(\tilde{W} - P_h^0(\tilde{W}), \tilde{V}_h - P_h^0(\tilde{W}))(t)
\leq K h_i^{\min\{s,k\}} \left[ \| (\tilde{V}_h - P_h^0(\tilde{V}))(t) \|_{L^2(\Omega)^3}^2 + \| (\tilde{\sigma}_h - P_h^2(\tilde{\sigma}))(t) \|_{L^2(\Omega)^9}^2 \right]^{1/2} \| (\tilde{V}, \tilde{\sigma})(t) \|_{PH^{s+1}(\Omega)^9}.
$$

**Proof** The two previous Lemmata 6 and 7 imply together the following estimates:

$$
\int_{\mathcal{S}_i} (\tilde{V}_h - P_h^1(\tilde{V})) |_{|r|_i} P_{il} (\tilde{\sigma} - P_h^2(\tilde{\sigma})) |_{|r|_i} ds \leq C \| P_{il} \| \ h_i^{\min\{s,k\}} \| \tilde{V}_h - P_h^1(\tilde{V}) \|_{L^2(\mathcal{T}_i)^3} \| \tilde{\sigma} \|_{H^{s+1}(\mathcal{T}_i)^6},
$$

$$
\int_{\mathcal{S}_i} (\tilde{\sigma}_h - P_h^2(\tilde{\sigma})) |_{|r|_i} P_{il} (\tilde{V} - P_h^1(\tilde{V})) |_{|r|_i} ds \leq C \| P_{il} \| \ h_i^{\min\{s,k\}} \| \tilde{\sigma}_h - P_h^2(\tilde{\sigma}) \|_{L^2(\mathcal{T}_i)^9} \| \tilde{V} \|_{H^{s+1}(\mathcal{T}_i)^3},
$$

$$
\int_{\mathcal{S}_i} (\tilde{V}_h - P_h^1(\tilde{V})) |_{|r|_i} B_{il} (\tilde{V} - P_h^1(\tilde{V})) |_{|r|_i} ds \leq C \ h_i^{\min\{s,k\}} \| \tilde{V}_h - P_h^1(\tilde{V}) \|_{L^2(\mathcal{T}_i)^3} \| \tilde{\sigma} \|_{H^{s+1}(\mathcal{T}_i)^6},
$$

$$
\int_{\mathcal{S}_i} (\tilde{\sigma}_h - P_h^2(\tilde{\sigma})) |_{|r|_i} B_{il} (\tilde{V} - P_h^1(\tilde{V})) |_{|r|_i} ds \leq C \ h_i^{\min\{s,k\}} \| \tilde{\sigma}_h - P_h^2(\tilde{\sigma}) \|_{L^2(\mathcal{T}_i)^9} \| \tilde{V} \|_{H^{s+1}(\mathcal{T}_i)^3}.
$$

where $C = C_1 C_2$, $l \in \mathcal{V}(i) \cup \mathcal{E}(i)$ in (103), $l \in \mathcal{V}(i) \cup \mathcal{K}(i) \cup \mathcal{E}(i)$ in (104), $l \in \mathcal{E}(i)$ in (105)-(106). In addition, applying (83), we have in the same way, for all $l \in \mathcal{V}(i)$:

$$
\int_{S_{il}} (\tilde{V}_h - P_h^1(\tilde{V})) |_{|r|_i} \beta_{il} (\tilde{\sigma} - P_h^2(\tilde{\sigma})) |_{|r|_i} ds \leq C \| \beta_{il} \| \ h_i^{\min\{s,k\}+1/2} \ h_i^{\min\{s,k\}} \| \tilde{V}_h - P_h^1(\tilde{V}) \|_{L^2(\mathcal{T}_i)^3} \| \tilde{\sigma} \|_{H^{s+1}(\mathcal{T}_i)^6},
$$

$$
\int_{S_{il}} (\tilde{\sigma}_h - P_h^2(\tilde{\sigma})) |_{|r|_i} \beta_{il} (\tilde{V} - P_h^1(\tilde{V})) |_{|r|_i} ds \leq C \| \beta_{il} \| \ h_i^{\min\{s,k\}+1/2} \ h_i^{\min\{s,k\}} \| \tilde{\sigma}_h - P_h^2(\tilde{\sigma}) \|_{L^2(\mathcal{T}_i)^9} \| \tilde{V} \|_{H^{s+1}(\mathcal{T}_i)^3}.\]

Thus, applying Theorem 2, we have a family of unstructured meshes satisfying (82) and (83).

Let $(\tilde{V}, \tilde{\sigma})$ be the exact solution of the symmetric pseudo-conservative problem (6)-(7)-(8)-(9) supposed to belong to $C^0([0,T], PH^{s+1}(\Omega)^9)$ for $s \geq 0$, and let $(\tilde{V}_h, \tilde{\sigma}_h) \in C^1([0,T], X^1_h \times X^2_h)$ be the solution of the semi-discrete problem (25)-(26)-(87). Then, there exists a constant $K > 0$ independent of $h$ such that

$$
\max_{t \in [0,T]} \left[ \| (\tilde{V} - \tilde{V}_h)(t) \|_{L^2(\Omega)^3}^2 + \| (\tilde{\sigma} - \tilde{\sigma}_h)(t) \|_{L^2(\Omega)^9}^2 \right]^{1/2} \leq K T h_i^{\min\{s,k\}} \| (\tilde{V}, \tilde{\sigma}) \|_{C^0([0,T], PH^{s+1}(\Omega)^9)}.
$$

(107)
Proof 10 First, applying Lemma 5 to the vector \( \tilde{W}_h - P_h^0 (\tilde{W}) = (\tilde{V}_h - P_h^1 (\tilde{V}), \tilde{\sigma}_h - P_h^2 (\tilde{\sigma})) \in X_h^1 \times X_h^2 \), we get

\[
a(\tilde{W}_h - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})) + b(\tilde{W}_h - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})) \geq 0,
\]

and, then

\[
m(\partial_t (\tilde{W}_h - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) \leq m(\partial_t (\tilde{W}_h - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) + a(\tilde{W}_h - \tilde{W}, \tilde{W}_h - P_h^0 (\tilde{W})) + b(\tilde{W}_h - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})).
\]

(108)

Thus, since \( \tilde{W}_h - P_h^0 (\tilde{W}) \in X_h^1 \times X_h^2 \), the difference between (91) and (92) gives us

\[
m(\partial_t (\tilde{W}_h - \tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W}))) + a(\tilde{W}_h - \tilde{W}, \tilde{W}_h - P_h^0 (\tilde{W})) + b(\tilde{W}_h - \tilde{W}, \tilde{W}_h - P_h^0 (\tilde{W})) = 0.
\]

After, we substract the previous equality in the right-hand side of (108) and we obtain

\[
m(\partial_t (\tilde{W}_h - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) \leq m(\partial_t (\tilde{W}_h - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) + a(\tilde{W}_h - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})) + b(\tilde{W}_h - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})).
\]

(109)

As \( P_h^1 (\tilde{V}) \) (resp. \( P_h^2 (\tilde{\sigma}) \)) is the orthogonal projection of \( \tilde{V} \) (resp. \( \tilde{\sigma} \)) onto \( X_h^1 \) (resp. \( X_h^2 \)) and \( \tilde{W}_h - P_h^0 (\tilde{W}) \in X_h^1 \times X_h^2 \), we necessarily have

\[
a(\tilde{W} - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W})) = 0
\]

and

\[
m(\partial_t (\tilde{W} - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) = \frac{1}{2} \frac{d}{dt} m((\tilde{W} - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) = 0.
\]

(111)

On the other hand, we can check that

\[
m(\partial_t (\tilde{W}_h - P_h^0 (\tilde{W})), \tilde{W}_h - P_h^0 (\tilde{W})) = \frac{1}{2} \sum_{i=1}^{N_T} \frac{d}{dt} \int_{t_i}^{t_{i+1}} \left[ \rho_i \left| \tilde{V}_h - P_h^1 (\tilde{V}) \right|^2 (t) + \left( \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) \Lambda_i \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) \right)^2 (t) \right] \, dt
\]

(112)

so, by taking also into account (110) and (111) in (109), we obtain

\[
\frac{d}{dt} \left[ \left( \left( \tilde{V}_h - P_h^1 (\tilde{V}) \right) (t) \right)^2_{L^2(\Omega)^3} + \left( \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) (t) \right) \right]_{L^2(\Omega)^6} \right] \, dt
\]

\[
\leq 2 \left( \min_{1 \leq i \leq N_T} \{ \rho_i, r_-(\Lambda_i) \} \right)^{-1} b(\tilde{W} - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W}))(t).
\]

(113)

Then, we integrate (113) from 0 to \( t \in [0, T] \) and, since (9) and (87) hold, we have

\[
\left( \left( \tilde{V}_h - P_h^1 (\tilde{V}) \right) (t) \right)^2_{L^2(\Omega)^3} + \left( \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) (t) \right) \right]_{L^2(\Omega)^6} \right] \, dt
\]

\[
\leq 2 \left( \min_{1 \leq i \leq N_T} \{ \rho_i, r_-(\Lambda_i) \} \right)^{-1} \int_0^T b(\tilde{W} - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W}))(t) \, dt.
\]

(114)

Therefore, the last inequality becomes

\[
\left( \left( \tilde{V}_h - P_h^1 (\tilde{V}) \right) (t) \right)^2_{L^2(\Omega)^3} + \left( \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) (t) \right) \right]_{L^2(\Omega)^6} \right] \, dt
\]

\[
\leq 2 \left( \min_{1 \leq i \leq N_T} \{ \rho_i, r_-(\Lambda_i) \} \right)^{-1} \int_0^T b(\tilde{W} - P_h^0 (\tilde{W}), \tilde{W}_h - P_h^0 (\tilde{W}))(t) \, dt
\]

(115)

which implies, by Lemma 8, that there exists \( K > 0 \) such that

\[
\max_{t \in [0, T]} \left[ \left( \left( \tilde{V}_h - P_h^1 (\tilde{V}) \right) (t) \right)^2_{L^2(\Omega)^3} + \left( \left( \tilde{\sigma}_h - P_h^2 (\tilde{\sigma}) \right) (t) \right) \right]_{L^2(\Omega)^6} \right]^{1/2} \leq K T h^{\min \{s, k\}} \left( \bar{V}, \tilde{\sigma} \right) \| \tilde{V} \|_{C^0([0, T], PH^{s+k}(\Omega)^3)},
\]

(117)
On the other hand, we also have by Lemma 6 the existence of $C > 0$ such that: $\forall t \in [0, T],$
\[
\left[ \left\| (\tilde{V} - P^1_h(\tilde{V}))(t) \right\|_{L^2(\Omega)^3}^2 + \left\| \tilde{\sigma} - P^1_h(\tilde{\sigma}))(t) \right\|_{L^2(\Omega)^6}^2 \right]^{1/2} \leq C h^{\min\{s,k\}+1} \left( \left\| \tilde{V}(t) \right\|_{PH^{s+1}(\Omega)^3}^2 + \left\| \tilde{\sigma}(t) \right\|_{PH^{s+1}(\Omega)^6}^2 \right)^{1/2},
\]
which combined to (117) concludes the proof.

4.3 Convergence of the totally discretized problem

In this section, we prove the convergence of the totally discretized scheme (42), following the method proposed in [56] for the second order leap-frog scheme.

**Theorem 3** Let $(\tilde{V}^n_h, \tilde{\sigma}^{n+1/2}_h)$ be the solution associated with (42) according to (11)-(12).

1. Let $(\tilde{V}_h, \tilde{\sigma}_h) \in C^{N+1}(0,T] \times X^1_h \times X^2_h)$ be the solution of the semi-discrete problem (25)-(26)-(87), and let $(\tilde{V}, \tilde{\sigma}) \in C^{N+1}(0,T] \cap C^0([0,T], PH^{s+1}(\Omega)^9)$ be the exact solution of (6)-(7)-(8)-(9). Under a condition on the time step $\Delta t$ of the same type as (75), there exists a constant $C > 0$ such that the following error estimate holds:

\[
\max_{n=0, \ldots, N} \left( \left\| \tilde{V}(t) - \tilde{V}^n_h \right\|_{C^0([0,T],L^2(\Omega)^3)}^2 + \left\| \tilde{\sigma}(t) - \tilde{\sigma}^n_h \right\|_{C^0([0,T],L^2(\Omega)^6)}^2 \right)^{1/2} \leq C \Delta t \left( \left\| \tilde{V}(t) - \tilde{V}^n_h \right\|_{C^0([0,T],L^2(\Omega)^3)} + \left\| \tilde{\sigma}(t) - \tilde{\sigma}^n_h \right\|_{C^0([0,T],L^2(\Omega)^6)} \right).
\]

2. Let $(\tilde{V}_h, \tilde{\sigma}_h) \in C^{\max\{\alpha,3\}}(0,T] \times X^1_h \times X^2_h)$ be the solution of the semi-discrete problem (25)-(26)-(87), where $\alpha \in \{1, 2, 3, 4\}$ is the coefficient defined at (37). Let $(\tilde{V}, \tilde{\sigma}) \in C^{\max\{\alpha,3\}}(0,T] \cap C^0([0,T], PH^{s+1}(\Omega)^9)$ be the exact solution of (6)-(7)-(8)-(9). Under a condition on the time step $\Delta t$ of the same type as (76), there exists a constant $C > 0$ such that the following error estimate holds:

\[
\max_{n=0, \ldots, N} \left( \left\| \tilde{V}(t) - \tilde{V}^n_h \right\|_{C^0([0,T],L^2(\Omega)^3)}^2 + \left\| \tilde{\sigma}(t) - \tilde{\sigma}^n_h \right\|_{C^0([0,T],L^2(\Omega)^6)}^2 \right)^{1/2} \leq C \Delta t \max\{\alpha,2\} \left( \left\| \tilde{V}(t) - \tilde{V}^n_h \right\|_{C^0([0,T],L^2(\Omega)^3)} + \left\| \tilde{\sigma}(t) - \tilde{\sigma}^n_h \right\|_{C^0([0,T],L^2(\Omega)^6)} \right).
\]

**Proof 11** Let be $\tilde{v}^n_i$ and $\tilde{\sigma}^{n+1/2}_i$ similarly defined as in (11)-(12). Then, for each tetrahedron $T_i$, we set $\tilde{v}^n_i = (\tilde{v}^n_i)_{1 \leq j \leq \text{dof}}$ and $\tilde{\sigma}^{n+1/2}_i = (\tilde{\sigma}^{n+1/2}_i)_{1 \leq j \leq \text{dof}}$, and we note the column vectors $\tilde{v}^n_i = (\tilde{v}^n_i)_{1 \leq i \leq N_T}$ and $\tilde{\sigma}^{n+1/2}_i = (\tilde{\sigma}^{n+1/2}_i)_{1 \leq i \leq \text{dof}}.$ We are interested in the estimation of the following local consistency error:

\[
\epsilon^n_i = \left( \| \tilde{v}^n_i - \tilde{V}_h(t) \|_{L^2(\Omega)^3}^2 + \| \tilde{\sigma}^{n+1/2}_i - \tilde{\sigma}_h(t) \|_{L^2(\Omega)^6}^2 \right)^{1/2},
\]

where $(\tilde{V}_h, \tilde{\sigma}_h) \in C^{N+1}(0,T] \times X^1_h \times X^2_h)$ is the semi-discrete solution of (25)-(26)-(87) and $(\tilde{v}^n_i, \tilde{\sigma}^{n+1/2}_i)$ has been computed as follows for all $N$ (even):

\[
\begin{align*}
M_p \frac{\tilde{v}^{n+1} - \tilde{V}(t)}{\Delta t} &= Q_N \tilde{\sigma}(t) - R_N \tilde{V}(t), \\
M_A \frac{\tilde{\sigma}^{n+1/2} - \tilde{\sigma}(t)}{\Delta t} &= X_N \tilde{V}(t) - Y_N \tilde{\sigma}(t).
\end{align*}
\]

Using Taylor formulae, there exist $c_{n+1/2} \in ]t_{n+1/2}, t_{n+1}[ \text{ and } d_{n+1/2} \in ]t_n, t_{n+1/2}[ \text{ such that }$

\[
\begin{align*}
\tilde{V}(t_{n+1}) &= \sum_{k=0}^{N} \frac{1}{k!} \left( \frac{\Delta t}{2} \right)^k \partial^k_t \tilde{V}(t_{n+1/2}) + \frac{1}{(N+1)!} \left( \frac{\Delta t}{2} \right)^{N+1} \partial^{N+1}_t \tilde{V}(c_{n+1/2}), \\
\tilde{V}(t_n) &= \sum_{k=0}^{N} \frac{(-1)^k}{k!} \left( \frac{\Delta t}{2} \right)^k \partial^k_t \tilde{V}(t_{n+1/2}) - \frac{1}{(N+1)!} \left( \frac{\Delta t}{2} \right)^{N+1} \partial^{N+1}_t \tilde{V}(d_{n+1/2}).
\end{align*}
\]
On the other hand, applying (123), (124) and (125) in the following equality
\[
M_\rho \left( \hat{\nabla}(t_{n+1}) - \hat{\nabla}(t_n) \right) = M_\rho \left( \hat{\nabla}(t_{n+1}) - \hat{\nabla}(t_n) \right) + M_\rho \left( \hat{\nabla}(t_n) - \hat{\nabla}^{n+1} \right) \\
= 2M_\rho \sum_{k=1,odd}^{N-1} \frac{1}{k!} \left( \frac{\Delta t}{2} \right)^k \partial_t^k \hat{\nabla}(t_{n+\frac{1}{2}}) - \Delta t Q_N \hat{\sigma}(t_{n+\frac{1}{2}}) + \Delta t R_N \hat{\nabla}(t_{n+\frac{1}{2}}) \\
+ \frac{1}{(N+1)!} \left( \frac{\Delta t}{2} \right)^{N+1} M_\rho \left( \partial_t^{N+1} \hat{\nabla}(c_{n+\frac{1}{2}}) + \partial_t^{N+1} \hat{\nabla}(d_{n+\frac{1}{2}}) \right). \tag{126}
\]

Using the identities (41) and (43), we deduce that the three first terms vanish. Then for each tetrahedron \( T_i \), it holds
\[
\hat{V}_i(t_{n+1}) - \hat{\nabla}^{n+1} = \sum_{j=1}^{dof} \left( \hat{V}_{ij}(t_{n+1}) - \hat{\nabla}^{n+1} \right) \hat{\sigma}_{ij} = \frac{1}{(N+1)!} \left( \frac{\Delta t}{2} \right)^{N+1} \left( \partial_t^{N+1} \hat{V}_i(c_{n+\frac{1}{2}}) + \partial_t^{N+1} \hat{V}_i(d_{n+\frac{1}{2}}) \right). \tag{127}
\]

Therefore, there exists a constant \( C > 0 \) such that:
\[
\| \hat{V}_h(t_{n+1}) - \hat{\nabla}^{n+1} \|_{L^2(\Omega)} \leq C \Delta t^{N+1} \| \hat{V}_h \|_{C^{N+1}(0,T],L^2(\Omega))^3}. \tag{128}
\]

In the very same way, we obtain analogous estimates for \( \| \hat{\sigma}_h(t_{n+\frac{1}{2}}) - \hat{\nabla}^{n+\frac{1}{2}} \|_{L^2(\Omega)} \). Finally, as \((\hat{V}_h, \hat{\sigma}_h)\) is a discrete approximation of \((\hat{V}, \hat{\sigma})\) by Th. 2, the estimation of the local consistency error is bounded in the following way:
\[
\hat{\varepsilon}_h^n \leq C \Delta t^{N+1} (\| \hat{V}_h, \hat{\sigma}_h \|_{C^{N+1}(0,T],L^2(\Omega))^p}). \tag{129}
\]

On the other hand, according to (129), we also have
\[
\begin{align*}
M_\rho \frac{\hat{V}(t_{n+1}) - \hat{V}(t_n)}{\Delta t} &= M_\rho \frac{\hat{V}(t_{n+1}) - \hat{\nabla}^{n+1}}{\Delta t} + Q_N \hat{\sigma}(t_{n+\frac{1}{2}}) - R_N \hat{\nabla}(t_{n+\frac{1}{2}}), \\
M_\Lambda \frac{\hat{\sigma}(t_{n+\frac{1}{2}}) - \hat{\sigma}(t_{n+\frac{1}{2}})}{\Delta t} &= M_\Lambda \frac{\hat{\sigma}(t_{n+\frac{1}{2}}) - \hat{\nabla}^{n+\frac{1}{2}}}{\Delta t} + X_N \hat{\nabla}(t_{n+1}) - Y_N \hat{\sigma}(t_{n+1}).
\end{align*} \tag{130}
\]

Then, using (42) and (130), the field \( (\hat{\nabla}^n, \hat{\nabla}^{n+\frac{1}{2}}) = (\hat{V}(t_n) - \hat{\nabla}^n, \hat{\sigma}(t_{n+\frac{1}{2}}) - \hat{\nabla}^{n+\frac{1}{2}}) \) satisfies the following set of equations:
\[
\begin{align*}
M_\rho \frac{\hat{\nabla}^{n+1} - \hat{\nabla}^n}{\Delta t} &= M_\rho \frac{\hat{\nabla}(t_{n+1}) - \hat{\nabla}^{n+1}}{\Delta t} + Q_N \hat{\nabla}^{n+\frac{1}{2}} - R_N \left[ \hat{\nabla}(t_{n+\frac{1}{2}}) - \hat{\nabla}^{n+\frac{1}{2}} \right], \\
M_\Lambda \frac{\hat{\nabla}^{n+\frac{1}{2}} - \hat{\nabla}^{n+\frac{1}{2}}}{\Delta t} &= M_\Lambda \frac{\hat{\sigma}(t_{n+\frac{1}{2}}) - \hat{\nabla}^{n+\frac{1}{2}}}{\Delta t} + X_N \hat{\nabla}^{n+1} - Y_N \left[ \hat{\sigma}(t_{n+1}) - \hat{\nabla}^{n+1} \right].
\end{align*} \tag{131}
\]

1. Firstly, we consider the case where \( N \) is even and the domain \( \Omega \) with free surface boundary conditions only. Then, \( R_N = 0 \) and \( Y_N = 0 \) in (131) by Remark 3, and we introduce the error energy \( \varepsilon^n \) at time \( n\Delta t \) by:
\[
\varepsilon^n = \frac{1}{2} \left( \hat{\nabla}^n \right)^t M_\rho \hat{\nabla}^n + \left( \hat{\nabla}^{n+\frac{1}{2}} \right)^t M_\Lambda \hat{\nabla}^{n+\frac{1}{2}}. \tag{132}
\]

Reasoning similarly as for the discrete energy \( E^n \) at (57), we can prove the following estimate under a condition of the same type as (75) and thanks to (129): there exists \( C > 0 \) such that
\[
\left( \| \hat{\nabla}^n \|_{L^2(\Omega)}^2 + \| \hat{\nabla}^{n+\frac{1}{2}} \|_{L^2(\Omega)}^2 \right)^{1/2} \leq C \Delta t^N (\| \hat{V}_h, \hat{\sigma}_h \|_{C^{N+1}(0,T],L^2(\Omega))^p}). \tag{133}
\]

Finally, we can conclude that the order of convergence of the scheme is \( \mathcal{O}(\Delta t^N + h_{\text{min}}(s,k)) \) by using the triangular inequality and results already established in the semi-discrete case (see Th. 2).
2. Then, for domains with absorbing boundary conditions (or mixed boundary conditions), we set \( N = 2 \). Therefore, it remains to evaluate the following term in (131) with (37) and \( \alpha \in \{1, 2, 3, 4\} \):

\[
\hat{\mathbf{v}}(t_{n+\frac{1}{2}}) - \hat{\mathbf{v}}[n+\frac{1}{2}] = \left( \hat{\mathbf{v}}(t_{n+\frac{1}{2}}) - \sum_{i=1}^{\alpha} b_i \hat{\mathbf{v}}(t_{n+1-i}) \right) + \hat{\mathbf{v}}[n+\frac{1}{2}].
\]  

(134)

For \((\hat{V}_h, \hat{\sigma}_h) \in C^{\max(\alpha, 3)}([0, T], X^1_h \times X^2_h)\), we know by a Taylor formula that there exists \( \xi_{i,n+\frac{1}{2}} \in ]t_{n+1-i}, t_{n+\frac{1}{2}}[ \), with \( i \in \{1, \ldots, \alpha\} \), such that

\[
\hat{\mathbf{v}}(t_{n+1-i}) = \sum_{k=0}^{\alpha-1} \frac{1}{k!} \left( \frac{1}{2} - i \right)^k \partial^k_t \hat{\mathbf{v}}(t_{n+\frac{1}{2}}) + \frac{1}{\alpha!} \left( \frac{1}{2} - i \right)^\alpha \Delta^\alpha_t \partial^\alpha_t \hat{\mathbf{v}}(\xi_{i,n+\frac{1}{2}}),
\]

(135)

which implies according to (38) that:

\[
\sum_{i=1}^{\alpha} b_i \hat{\mathbf{v}}(t_{n+1-i}) = \hat{\mathbf{v}}(t_{n+\frac{1}{2}}) + \frac{\Delta t^\alpha}{\alpha!} \left[ \sum_{i=1}^{\alpha} b_i \left( \frac{1}{2} - i \right)^\alpha \partial^\alpha_t \hat{\mathbf{v}}(\xi_{i,n+\frac{1}{2}}) \right].
\]

(136)

And, finally, in a similar way for the term \( \hat{\sigma}(t_{n+1}) - \hat{\sigma}^{n+1} \), the system (131) can be rewritten as

\[
M_p \frac{\hat{\mathbf{v}}^{n+1} - \hat{\mathbf{v}}^n}{\Delta t} = Q \hat{\mathbf{v}}^{n+\frac{1}{2}} - R \hat{\mathbf{v}}^{[n+\frac{1}{2}]} + M_p \frac{\hat{\mathbf{v}}(t_{n+1}) - \hat{\mathbf{v}}^{n+1}}{\Delta t} + \frac{\Delta t^\alpha}{\alpha!} R \left[ \sum_{i=1}^{\alpha} b_i \left( \frac{1}{2} - i \right)^\alpha \partial^\alpha_t \hat{\mathbf{v}}(\xi_{i,n+\frac{1}{2}}) \right],
\]

(137)

\[
M_A \frac{\hat{\sigma}^{n+\frac{1}{2}} - \hat{\sigma}^{n+\frac{1}{2}}}{\Delta t} = X \hat{\mathbf{v}}^{n+1} - Y \hat{\mathbf{v}}^{n+\frac{1}{2}} + M_A \frac{\hat{\sigma}(t_{n+\frac{1}{2}}) - \hat{\sigma}^{n+\frac{1}{2}}}{\Delta t} + \frac{\Delta t^\alpha}{\alpha!} Y \left[ \sum_{i=1}^{\alpha} b_i \left( \frac{1}{2} - i \right)^\alpha \partial^\alpha_t \hat{\sigma}(\xi_{i,n+1}) \right].
\]

(138)

Now, when \( \alpha = 1 \), we introduce the error energy \( \mathcal{E}^n \) at time \( n\Delta t \) by:

\[
\mathcal{E}^n = \frac{1}{2} \left[ (\hat{\mathbf{v}}^n)^t M_p \hat{\mathbf{v}}^n + (\hat{\mathbf{v}}^{n+\frac{1}{2}})^t M_A \hat{\mathbf{v}}^{n+\frac{1}{2}} \right] - \frac{\Delta t}{4} \left[ (\hat{\mathbf{v}}^n)^t R \hat{\mathbf{v}}^n - (\hat{\mathbf{v}}^{n+\frac{1}{2}})^t Y \hat{\mathbf{v}}^{n+\frac{1}{2}} \right].
\]

(139)

elsewhere, for \( \alpha \in \{2, 3, 4\} \), \( \mathcal{E}^n \) is defined by:

\[
\mathcal{E}^n = \frac{1}{2} \left[ (\hat{\mathbf{v}}^n)^t M_p \hat{\mathbf{v}}^n + (\hat{\mathbf{v}}^{n+\frac{1}{2}})^t M_A \hat{\mathbf{v}}^{n+\frac{1}{2}} \right] - \frac{\Delta t}{4} \sum_{i=1}^{\alpha} |b_i| \left[ (\hat{\mathbf{v}}^n)^t R \hat{\mathbf{v}}^n - (\hat{\mathbf{v}}^{n+\frac{1}{2}})^t Y \hat{\mathbf{v}}^{n+\frac{1}{2}} \right] + \frac{\Delta t}{4} \sum_{i=2}^{\alpha} |b_i| \left( \hat{\mathbf{v}}^n - \gamma_i \hat{\mathbf{v}}^{n+1-i} \right)^t R \left( \hat{\mathbf{v}}^n - \gamma_i \hat{\mathbf{v}}^{n+1-i} \right).
\]

(140)

Reasoning similarly as for the discrete corrected energy \( \mathcal{E}^n \) at (63) or (64), we can prove the following estimate under a condition of the same type as (76) and thanks to (129): there exists \( C > 0 \) such that

\[
\left( \left\| \hat{\mathbf{v}}^n \right\|_{L^2(\Omega)}^2 + \left\| \hat{\mathbf{v}}^{n+\frac{1}{2}} \right\|_{L^2(\Omega)}^2 \right)^{1/2} \leq C \left( \Delta t^\alpha \left\| (\hat{V}, \hat{\sigma}) \right\|_{C^\alpha([0, T], L^2(\Omega)^2)} + \Delta t^2 \left\| (\hat{V}, \hat{\sigma}) \right\|_{C^1([0, T], L^2(\Omega)^2)} \right) .
\]

(141)

Finally, we can conclude that the order of convergence of the scheme is \( \mathcal{O}(\Delta t^\min(\alpha, 2) + h^\min(\alpha, k)) \) by using the triangular inequality and results already established in the semi-discrete case (see Th. 2).

\[\square\]
5 Numerical results

The DG method has been applied to several numerical problems, each to illustrate a specific property. First, the propagation of an eigenmode aims to realize a convergence study of the schemes and a comparison with the theoretical analysis. Moreover, it allows validating the free surface conditions and comparing the properties of the standard and fourth-order leap-frog schemes. The second study concerns the propagation of a pulse and focuses on the absorbing boundary conditions. The last application concerns the propagation of the wave produced by an explosive source in a half space which enables, firstly, validating the source introduction and apply the method in a more realistic context.

5.1 Convergence study. Propagation of an eigenmode in 3D

The first problem concerns the propagation of an eigenmode in three dimensions of space. The domain of computation is the unit cubic cavity on which we apply free surface boundary conditions. We are interested in the (1, 1, 1) mode whose exact solution at time \( t \) and in \( X = (x, y, z) \) is given by:

\[
\begin{align*}
V_x(t, X) &= \cos(\pi x) \left[ \sin(\pi y) - \sin(\pi z) \right] \cos(\Omega t), \\
V_y(t, X) &= \cos(\pi y) \left[ \sin(\pi z) - \sin(\pi x) \right] \cos(\Omega t), \\
V_z(t, X) &= \cos(\pi z) \left[ \sin(\pi y) - \sin(\pi x) \right] \cos(\Omega t), \\
\sigma_{xx}(t, X) &= -A \sin(\pi x) \left[ \sin(\pi y) - \sin(\pi z) \right] \sin(\Omega t), \\
\sigma_{yy}(t, X) &= -A \sin(\pi y) \left[ \sin(\pi z) - \sin(\pi x) \right] \sin(\Omega t), \\
\sigma_{zz}(t, X) &= -A \sin(\pi z) \left[ \sin(\pi x) - \sin(\pi y) \right] \sin(\Omega t), \\
\sigma_{xy}(t, X) &= \sigma_{xz}(t, X) = \sigma_{yz}(t, X) = 0,
\end{align*}
\]

with \( A = \sqrt{2/\rho \mu} \) and \( \Omega = \pi \sqrt{2/\rho \mu} \). We set dimensionless values for the medium properties, \( \rho = 1.0 \), \( \lambda = 0.5 \) and \( \mu = 0.25 \), which implies that \( v_p=1 \) and \( v_s=0.5 \). The initialisation of the leap-frog scheme is realized by deducing, from the analytical expressions (142), values for the velocity components at \( t = 0 \) and for stress components at \( n = 0 \), where \( n \) is the number of time steps.

Series of different uniform meshes are constructed by dividing the domain in cubic cells which are split in six tetraedra. The mesh spacing \( h \) is the length of the longest edge of the mesh. The characteristics of all meshes can be found in Table 1. Different methods have been used, the notation \( P_p \) referring to a discretization based on polynomials of degree \( p \). Two leap-frog schemes are also compared: the standard one, second-order accurate (referred as LF2), and a fourth-order extension (named LF4 in the following). Simulations are performed until \( t_{max} = 142 \) s which corresponds to 50 periods in the cube and is a sufficiently long time to distinguish the results from different time and space discretizations.

For all the applications, the time-step of the simulation is the minimum of all local time-steps calculated in each tetraedron using the following formula

\[
\Delta t_{LF2} = \min_{T_i} \Delta t_i \quad \text{with} \quad \Delta t_i = \frac{h_i}{2p + 1 \cdot v_p |_{T_i}} \quad \text{and} \quad \Delta t_{LF4} = 2.5 \times \Delta t_{LF2},
\]

which depends on the space approximation \( p \), \( h_i \), the smallest height in the tetraedron and the local value of \( v_p \). The ratio between \( \Delta t_{LF1} \) and \( \Delta t_{LF2} \) has been deduced experimentally in [35].

We present in Figure 1 the results of the convergence study. First, Figures (a) and (b) display the maximum value of the numerical error \( err_{L2} \), at time \( t_{max} \), and in logarithmic scale by respect to the mesh spacing \( h \) for LF2 and LF4 leap-frog schemes respectively. The corresponding results of efficiency, i.e. the error as a function of the CPU time, are presented in Figures (c) and (d). The values of the error, the corresponding convergence orders and the CPU times are specified in Table 2 and complete these figures. As can be noticed, some simulations are
Table 1: Characteristics of the uniform meshes used for the convergence study; \( N_N \) and \( N_T \) are respectively the total number of nodes and tetraedra of the meshes.

missing, those corresponding to the highest-degree space approximations and the finest meshes, simply for reason of CPU cost (these calculations have been done on a standard PC). In contrast, the solutions of the \( P_1 \) method show significant dispersion errors for the coarsest meshes in long time (the solution can be shifted more than over a period) and have been excluded for this reason.

Figure 1: Convergence study for the propagation of an eigenmode. Error as a function of the mesh size \( h \) (a and b) and CPU time (c and d) for different \( P_p-LF_i \) methods (\( p = 1, \ldots, 4, \ i = 2, 4 \)) and leap-frog time schemes (left column for LF2 and right column for LF4 respectively). Errors at \( t_{\text{max}} = 142 \) s using uniform meshes.

When examining the Figures 1, (a) and (b), we notice that the results are in accordance with the theoretical analysis. For a method \( P_p-LF_i \), in presence of free-surface conditions only, the error is \( O(h^p + \Delta t^i) \). Then, for \( p \leq 2 \), the global convergence rate is about \( p \) and dominated by the one in space, regardless of the leap-frog scheme; this is what is observed in both figures and confirmed by the values in Table 2. For the standard leap-frog scheme (LF2) and higher values of \( p \), the convergence rate is still equal to 2 because dominated by the one in time. The improvement brought by the LF4 leap-frog extension is clear since, for \( p \geq 3 \), the convergence rate is no more constrained by the time scheme and is also equal to \( p \), as noticeable in Figure (b) as well as in convergence orders of Table 2. Note that this study also validates the implementation of the free surface condition.

We are now interested in the efficiency of the different methods presented in Figures 1, (c) and (d). We remark that, to reach a given error level, methods based on high order polynomials generally need lower CPU times, for
both types of leap-frog schemes. However, when using the standard leap-frog scheme, the \( P_3 \) method becomes too expensive compared to the \( P_3 \) scheme, especially for fine meshes, mainly because of the limitation in convergence due to the time scheme and to a larger number of degrees of freedom. The improvement brought by the high-order leap-frog (LF4) extension is obvious and can be explained by better convergence rates for \( p \geq 3 \) and, also, by the use of greater time-steps. So, even if it is based on a multi-step algorithm, the additional iterations are counterbalanced by a greater time-step which makes the LF4 method more efficient. This is confirmed by the CPU times given in Table 2 proving a reduction of about 25% of the CPU time between results corresponding to LF2 and LF4 methods, even for \( p = 1, 2 \) for which error levels are comparable for LF2 and LF4 time discretizations.

In order to complete the comparison between both LF2 and LF4 leap-frog schemes, we plot in Figure 2 the error until \( t_{\text{max}} \), obtained using the mesh M4, for various \( P_p \) methods \((2 \leq p \leq 4)\). The Figures (a) and (b) corresponding to \( p = 2 \) and 3, justify the choice of a long time for simulations. For \( p = 2 \), the difference due to the time approximation is very low and only visible after \( 30 \) s. The difference between time methods is more visible for the \( P_3 \) method, in Figure (b); for times greater than \( 10 \) s, the error associated to the standard leap-frog scheme increases continuously since it remains quite constant with the LF4 extension. The improvement due to the LF4 leap-frog scheme is clear in Figures (c) and (d), corresponding to the \( P_4 \) method. The high-order leap-frog extension allows a reduction of the error by a factor of about 100 and proves the interest of using high-order space approximations.

Finally, to conclude this study, we plot in Figure 3, the value of the velocity \( V_x \) recorded inside the domain, at location \((0.25, 0.25, 0.25)\), over a long time, corresponding to 200 and 400 periods approximately. We compare the exact solution (series of dots) to the results obtained using the LF2 and LF4 leap-frog schemes, with the coarse mesh M2 and the \( P_1 \) method. For the first interval, Figure (a), the results are very close. We can notice a slight delay of the LF2 solution compared to the two others. For the second interval, Figure (b), which corresponds to a very long time, the difference between the solutions is clear. The dispersion error of the LF2 time scheme results in a much larger delay. Note that the amplitude is correct for both solutions and that the solution corresponding to the LF4 method is in perfect accordance with the exact one.

### 5.2 Propagation of a pulse. Comparison between two types of absorbing boundary conditions

The second problem we have studied is the propagation of a pulse in a homogeneous medium. Although this is a one-dimensional problem, it is possible to solve it in two or three dimensions of space. This is done here, in 3D, simply to apply the solver proposed in this paper. This study follows previous results presented by Delcourte et

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**Table 2:** Errors in \( L^2 \)-norm, CPU times and convergence orders corresponding to the convergence study, \( P_p \)-LF\( i \) methods \((p = 1, \ldots, 4, i = 2, 4)\).
al. [25], using a two-dimensional method, that highlighted some limitations of the standard absorbing boundary conditions and aims to focus on the improvement proposed in section 2.4.2. More precisely, we compare two absorbing boundary conditions among those defined by equations (39): one corresponding to $\alpha = 1$ and referred as “standard absorbing condition” and a second one ($\alpha = 2$), based on two levels in time and called “modified absorbing condition” in the following.

The domain of computation is $[0; 2] \times [0; l]^2$; the length $l$ being variable in y- and z-directions since uniform meshes are constructed always considering 5 nodes in these two directions, whatever the discretization chosen in the x-direction. As previously, dimensionless values are set for $\rho = 1$, $\lambda = 0.5$ and $\mu = 0.25$ resulting in $v_p = 1$. 

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Figure 2: Errors in $L^2$-norm as a function of time for LF2 and LF4 leap-frog schemes. (a) $P_2$ methods, (b) $P_3$ methods, (c) $P_4$ methods and (d) detail of figure (c) for $t \in [0; 40]$. All results correspond to mesh M4.

Figure 3: Velocity $V_x$ recorded at location $(0.25, 0.5, 0.25)$ as a function of time, for time intervals $[563.5; 566.5]$ (a) and $[1132.; 1135.]$ (b). $P_4$ method, mesh M2 and comparison between LF2 (dashed lines) and LF4 (solid lines) leap-frog schemes. 

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and \( v_s = 0.5 \). The compression pulse is placed in \( x_0 = 1.5 \) and is defined by

\[
V_x(x, t_0 = 0) = \exp^{-50(x-x_0)^2}, \quad \sigma_{xx}(x, t_0 = 0) = -\exp^{-50(x-x_0)^2},
\]

and, for all \( t \), the other components of velocity or stress are equal to 0. The analytical solution is calculated using the characteristic method according to the initial condition at \( t_0 \) and expresses

\[
\begin{align*}
V_x(x, t) &= \frac{1}{2} \left[ V_x(x - v_p t, 0) - \frac{\sigma_{xx}(x - v_p t, 0)}{\rho v_p} \right] + \frac{1}{2} \left[ V_x(x + v_p t, 0) + \frac{\sigma_{xx}(x + v_p t, 0)}{\rho v_p} \right], \\
\sigma_{xx}(x, t) &= \frac{1}{2} \rho v_p \left[ V_x(x + v_p t, 0) + \frac{\sigma_{xx}(x + v_p t, 0)}{\rho v_p} \right] - \frac{1}{2} \rho v_p \left[ V_x(x - v_p t, 0) - \frac{\sigma_{xx}(x - v_p t, 0)}{\rho v_p} \right].
\end{align*}
\]

The pulse propagates at \( v_p \) and the exact solution for velocity \( V_x \) at different times is shown in Figure 4. We compare the solutions obtained using different \( \mathcal{P}_p \) methods (\( p = 1, \ldots, 4 \)) and both standard and modified absorbing boundary conditions. In this study, we use the standard leap-frog time scheme in order to focus on the absorbing boundary conditions. We present, in Figure 5, the \( L^2 \)-error as a function of time for the various methods. These results correspond to a uniform mesh containing 2 025 nodes and 7 680 tetraedra (\( h_{\text{min}} = 2.5 \times 10^{-2}, h_{\text{max}} = 4.33 \times 10^{-2} \) and \( l = 1.0 \times 10^{-1} \)). When examining the different figures, we first notice that all results have the same behavior that can be related to the evolution of the pulse during time, as shown in Figure 4. The error increases rapidly and stabilizes around the value \( t = 1 \) s that corresponds to times for which the pulse propagates inside the computational domain. Note the vertical scale of each figure which indicates a global reduction of the error for higher values of \( p \). For \( t > 1 \) s, the error decreases as the pulse begins to go out of the domain. Until \( t = 1.5 \) s, the error only depends on the accuracy of the method (i.e. the value of \( p \)) and not on the absorbing boundary condition. For times greater than 1.5 s, the results corresponding to the two boundary conditions can be distinguished; the improvement brought by the modified absorbing condition is clearly visible. For \( p \geq 2 \), the error at \( t = 2 \) s is significantly lower with the modified absorbing condition which proves less spurious reflections than with the standard absorbing condition. It should also be noted that, when using the standard condition, the levels of error at \( t = 2 \) s corresponding to \( 2 \leq p \leq 4 \) are very close which proves that this condition tends to reduce the convergence.

For a detailed comparison of the solutions of the two absorbing methods, we plot in Figure 6 the value of \( V_x \) at \( t = 2s \), as a function of \( x \) in the middle of the domain, i.e. in \( y = z = l/2 \). We compare the exact solution (a series of dots) with the results of various \( \mathcal{P}_p \) methods (\( p = 1, \ldots, 4 \)) using the standard (dashed lines) and the modified (solid lines) absorbing conditions. At this time, the pulse is out of the computational domain and the various curves are a way to measure the level of accuracy of the boundary conditions. We present, in Figure 6 (a), the results of the \( \mathcal{P}_1 \) method: many spurious oscillations are visible and the two absorbing conditions lead to comparable error levels (\( \max_x |V_x| \simeq 4.4 \times 10^{-4} \)). Better results are obtained using the \( \mathcal{P}_2 \) method, in Figure 6 (b): the use of the modified absorbing condition reduces the amplitude of the oscillations (\( \max_x |V_x| \simeq 2.5 \times 10^{-4} \) with the standard and \( \max_x |V_x| \simeq 9.1 \times 10^{-5} \) with the modified absorbing condition) but a distinct reflection is visible at about \( x = 0.15 \) for both types of conditions. The Figures 6 (c) and (d) compare the solutions for \( \mathcal{P}_3 \) and \( \mathcal{P}_4 \) methods. In both cases, reflections are visible with the standard absorbing condition even if the amplitude is slightly reduced for the \( \mathcal{P}_3 \) method (\( \max_x |V_x| \simeq 2.0 \times 10^{-4} \) for \( \mathcal{P}_3 \) method and \( \max_x |V_x| \simeq 1.4 \times 10^{-4} \) for \( \mathcal{P}_4 \) method). The improvement brought by the modified absorbing condition is obvious and, at this scale, it is not
Figure 5: Errors in $L^2$-norm as a function of time for (a) $P_1$, (b) $P_2$, (c) $P_3$ and (d) $P_4$ methods. All results correspond to a uniform mesh of 2 025 nodes and 7 680 tetraedra. Dashed (respectively solid) lines for the standard (respectively modified) absorbing boundary condition.

possible to distinguish them from the exact solution. We plot in Figure 6 (c) the results obtained using the $P_3$ and $P_4$ methods and the modified absorbing condition. Some spurious reflections remain but with amplitudes clearly lower ($\max_x |V_x| \simeq 6.8 \times 10^{-6}$ for $P_3$ method and $\max_x |V_x| \simeq 3.8 \times 10^{-6}$ for $P_4$ method). This illustrates the improvement brought by for the absorbing boundary conditions proposed in section 2.4.2. The comparison of the two formulations, corresponding to $\alpha = 1$ and $\alpha = 2$, shows a clear reduction of the spurious reflections at the boundary in a simple case where the wave propagates perpendicularly to the boundary. This condition has to be studied in a more complex configuration, in particular involving razing waves. However, note that the modified absorbing boundary condition is very simple to implement as it requires the storage of some unknowns of the previous time level, only at the boundary.

5.3 Explosive source in a half space

The method has been applied to a last test case i.e. the propagation of the wave produced by an explosive source in a homogeneous half-space, as illustrated in figure 7. This case, which is really 3D, allows a validation of the source implementation and an additional study of the absorbing boundary conditions. The medium is homogeneous of density $\rho=2800$ kg/m$^3$ and P- and S-waves velocities are respectively $v_p=6400$ m/s and $v_s=3700$ m/s. We apply a free surface condition on the upper boundary of the domain and absorbing conditions on the other boundaries. The source is located at the center of the domain at 1000m depth (at the coordinates (0;0;1000)) and four surface stations are placed at 6000m from the epicenter, symmetrically on the $x$- and $y$-axis as described in figure 7. The source signal is a Ricker wavelet whose moment-rate is given by

$$s(t) = \frac{3.1 \times 10^{16}}{\alpha \sqrt{\pi}} \exp \left[ -\frac{(t-1)^2}{\alpha^2} \right]$$

with $\alpha=0.25$; it corresponds to a central frequency equal to $f_c=1.0$Hz and a maximum frequency of $f_{max}=3.0$Hz. Thus, the wavelength $L = v_s/f_{max}$ is approximatively equal to 1233m. The explosive source is introduced as a
right hand side on the diagonal components of the stress tensor i.e. for $\sigma_{xx}$, $\sigma_{yy}$ and $\sigma_{zz}$. In order to compare the results obtained using several methods $P_p$, especially the finite volume method ($P_0$), this right hand side writes

$$s(t) g(x,y,z)$$

where $g$ is defined by

$$g(x,y,z) = \frac{1}{M} \exp \left[ -\frac{(x-x_S)^2 + (y-y_S)^2 + (z-z_S)^2}{h^2} \right]$$

with $M = \int_{\Omega} g(x,y,z) \, dv$, where $(x_S, y_S, z_S)$ is the source coordinates. Note that the support of $g$ is taken sufficiently small, compared to the element size, for an accurate approximation of a point source, a right hand side based on a Dirac function $s(t) \delta(x_S, y_S, z_S)$ being only accurate for high degree approximation based methods. Initial conditions for the system are $\vec{V} = \vec{0}$ and $\sigma = \vec{0}$ and the solutions are calculated until $t=6.0$ s. We also point out that the finite-volume method presented in this study is a very standard one where the unknowns are set constant in the mesh elements and cannot be compared to high-order finite-volume methods, in particular the one presented by Dumbser et al.
The calculation domain is 36km × 36km × 12km. Solutions have been obtained using several uniform tetrahedral meshes obtained, as previously, by dividing the domain in cubic cells of edge \( h \) which are split in six tetraedra. The mesh spacing \( h \), i.e. here the smallest edge of the mesh, can be expressed as a function of the wavelength \( L \) as summarized in Table 3.

<table>
<thead>
<tr>
<th>( h ) (meters)</th>
<th>500</th>
<th>400</th>
<th>300</th>
<th>200</th>
<th>150</th>
<th>100</th>
<th>75</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \simeq L/n )</td>
<td>( L/2.5 )</td>
<td>( L/3 )</td>
<td>( L/4 )</td>
<td>( L/6 )</td>
<td>( L/8 )</td>
<td>( L/12 )</td>
<td>( L/16 )</td>
</tr>
<tr>
<td>( N_T )</td>
<td>746 496</td>
<td>1 458 000</td>
<td>3 456 000</td>
<td>11 664 000</td>
<td>27 648 000</td>
<td>93 312 000</td>
<td>221 184 000</td>
</tr>
<tr>
<td>( N_T/\text{proc.} )</td>
<td>11 664</td>
<td>22 781</td>
<td>54 000</td>
<td>182 250</td>
<td>432 000</td>
<td>1 458 000</td>
<td>3 456 000</td>
</tr>
</tbody>
</table>

Table 3: Correspondance between the mesh spacing \( h \) and the wavelength \( L \) of the source function. Characteristics of the different meshes (\( N_T \) total number of tetraedra, \( N_T/\text{proc.} \), number of tetraedra for each processor).

Calculations have been performed on 64 processors using message passing interface (MPI), whatever the mesh spacing, for a better comparison of the CPU times between the different methods. The characteristics of the different meshes are also listed in Table 3.

For this test, we have proposed some criteria in order to estimate the accuracy of the solution. Firstly, we check two particular properties, consequences of an explosive source, which are i) \( V_z = 0 \) at stations \( C_1 \) and \( C_3 \) (on y-axis) and \( V_y = 0 \) at stations \( C_2 \) and \( C_4 \) (on x-axis) and ii) the evolution at the surface of the vertical velocity component \( V_z \) only depends on the distance to the epicenter, thus the \( V_z \) profiles obtained at the four stations are superimpose. Moreover, for a validation of the source introduction, the amplitudes and the arrival times of the profiles are compared to a reference solution computed at these four stations (and not everywhere in the computational domain) using the Discrete Wave Number method [10]. Then, three different errors are defined. First, \( \text{err}_1 \) is concerned with the components which should ideally be equal to zero and expresses

\[
\text{err}_1 = \max_{t \in [0;6]} \left[ |V_x|_{C_1}, |V_y|_{C_2}, |V_z|_{C_3}, |V_y|_{C_4} \right].
\]

The two other errors depend on the values at the extrema, as illustrated in Figure 8. For each non null velocity profile (for instance, here, \( V_z \) for the reference solution at station \( C_1 \)), we identify three extrema and note their values \( (V_z(\text{extr}_i), i=1,2,3) \) and the corresponding times \( (t_i, i=1,2,3) \). From these data, we define a relative error, in the amplitude or in time, between two profiles (for instance profiles at stations \( C_j \) and \( C_k \)), for a given velocity component \( V_a \) and a given extremum \( \text{extr}_i \) by the following formulae, for amplitudes and times

\[
\text{err}_{rel}(V_a, C_j/C_k, \text{extr}_i)_{ampl} = \frac{|V_a(\text{extr}_i)_{C_j} - V_a(\text{extr}_i)_{C_k}|}{V_a(\text{extr}_i)_{ref}} \quad \text{err}_{rel}(V_a, C_j/C_k, \text{extr}_i)_t = \frac{t_i/C_j - t_i/C_k}{t_i/\text{ref}},
\]

where the subscript \( V_{ref} \) stands for the corresponding value of the reference solution. Thus, error \( \text{err}_2 \) measures the asymmetries of the computed solution (i.e. the mean maximum error of the amplitude and the travel time between solutions computed at two stations) of the vertical velocity component \( V_z \) which writes

\[
\text{err}_{2/ampl} = \frac{1}{3} \sum_{i=1}^{3} \max_{1 \leq j,k \leq 4} \text{err}_{rel}(V_z, C_j/C_k, \text{extr}_i)_{ampl}, \quad \text{err}_{2/t} = \frac{1}{3} \sum_{i=1}^{3} \max_{1 \leq j,k \leq 4} \text{err}_{rel}(V_z, C_j/C_k, \text{extr}_i)_t.
\]
Finally, \( err_3 \) is the mean error between the computed and the reference solutions at all stations, for all non null velocity components and for all extrema. This error writes, for instance for amplitudes

\[
err_{3/ampl} = \frac{1}{24} \sum_{i=1}^{3} \left[ \sum_{j \in \{2,4\}} err_{rel} (V_z, C_j/C_{j/ref}, extr_i)_{ampl} + \right.
\]

\[
\left. \sum_{j \in \{1,3\}} err_{rel} (V_y, C_j/C_{j/ref}, extr_i)_{ampl} + \sum_{j=1}^{4} err_{rel} (V_z, C_j/C_{j/ref}, extr_i) \right],
\]

an equivalent definition being applied to define \( err_{3/t} \).

Calculations have been done using methods based on different degrees of interpolation \((0 \leq p \leq 4)\) in combination with the meshes defined in table 3 in order to study the convergence and CPU time corresponding to the three types of errors; the results of these calculations are gathered in figure 9 and the CPU times of each calculation are given in table 4. Note that, as previously, \( \mathcal{P}_p \) refers to a method based on polynomials of degree \( p \) and the standard leap-frog time method. This first series of results is obtained using the standard absorbing boundary condition.

The first line of the figure presents the values of \( err_1 \) as a function of the mesh spacing \( h \) (figure a) and by respect to the corresponding CPU time (figure b). When examining these results, we notice (figure a) that \( err_1 \) corresponding to \( \mathcal{P}_0 \) method remains nearly constant, whatever the mesh spacing, at about \( 10\% \) of the maximum value of the velocity. Results are clearly improved by the use of higher degree methods since, in comparison with the \( \mathcal{P}_0 \) results, the value of \( err_1 \) of the \( \mathcal{P}_2 \) method is reduced by a factor \( 10^9 \). Moreover, higher degree methods \( \mathcal{P}_3 \) and \( \mathcal{P}_4 \) are also more efficient since a given error level of accuracy is obtained for lower CPU times.

Now, we study the second line of figures corresponding to \( err_2 \), which is a relative error traducing the asymmetry of the \( V_z \) profiles by respect to \( h \) (figure c) and the CPU time (figure d). As previously, we notice that the results of the \( \mathcal{P}_0 \) method are not accurate enough even using the finest meshes. The improvement on the symmetry is obvious when using the other methods, especially the methods with \( p \geq 2 \) resulting in \( err_2 \) values lower than \( 0.1\% \). As previously for \( err_1 \), \( \mathcal{P}_3 \) and \( \mathcal{P}_4 \) are the two most efficient methods.

Finally, we analyze the results obtained for \( err_3 \) (figures 9, e and f) which is a relative error between the computed and the reference solutions at the extrema. As previously, the lowest error levels correspond to the \( \mathcal{P}_2 \) to \( \mathcal{P}_4 \) methods whereas the \( \mathcal{P}_0 \) method produces relative errors in the order of \( 10\% \) and this independently on \( h \). The other methods allow errors lower than \( 5\% \) when using the finest meshes. Note that, for this last error, if the improvement between \( \mathcal{P}_0 \) and \( \mathcal{P}_1 \) is obvious, it is more limited for the highest degree methods. Furthermore, these methods lead to comparable levels of error. This fact is also visible when \( err_3 \) is plotted as a function of the CPU time (figure f). The reduction of the error levels for the highest values of \( p \) is not sufficient to compensate their overcost; thus, for this last criterion, the \( \mathcal{P}_1 \) method appears to be most efficient.

We also present, in figure 10, the errors \( err_{2/t} \) (figure a) and \( err_{3/t} \) (figure b) on arrival times as a function of the mesh spacing. This error is plotted only for \( \mathcal{P}_0 \) and \( \mathcal{P}_1 \) methods because it is equal to zero for \( 2 \leq p \leq 4 \), whatever the value of \( h \). Even for fine meshes, the arrival times are not exactly the same for \( p = 0 \) and \( p = 1 \). On figure b, which corresponds to \( err_{3/t} \), we clearly distinguish the \( \mathcal{P}_0 \) method for which the relative error is greater than \( 1\% \). This error is reduced using the \( \mathcal{P}_1 \) method and very low for \( p \geq 2 \). In first conclusion of this study, we have to distinguish results obtained with the \( \mathcal{P}_0 \) method from those of the other methods. First of all, the \( \mathcal{P}_0 \)
Figure 9: Convergence and CPU time for errors on amplitude, $err_1$ (a and b), $err_2$ (c and d) and $err_3$ (e and f). Explosive source in a half space, solutions at $t=6\,\text{s}$, uniform meshes, computations on 64 processors.

Figure 10: Convergence on arrival times, $err_2/t$ (figure a) and $err_3/t$ (figure b). Explosive source in a half space, solutions at $t=6\,\text{s}$, uniform meshes, computations on 64 processors.
Table 4: CPU times of the calculations for the $P_p$ methods and the different meshes on 64 processors.

<table>
<thead>
<tr>
<th>$P_p/h$</th>
<th>500</th>
<th>400</th>
<th>300</th>
<th>200</th>
<th>150</th>
<th>100</th>
<th>75</th>
</tr>
</thead>
<tbody>
<tr>
<td>$P_0$</td>
<td></td>
<td></td>
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<tr>
<td>$P_1$</td>
<td></td>
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<tr>
<td>$P_2$</td>
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<td></td>
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</tr>
<tr>
<td>$P_3$</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>$P_4$</td>
<td></td>
<td></td>
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</tr>
</tbody>
</table>

method leads to solutions that do not verify the symmetry condition, for amplitudes and arrival times, and the use of fine meshes do not allow improving these results. Better solutions are obtained with the $P_1$ method but the use of higher-degree space interpolations clearly helps improving the different types of errors. When combining the different criteria on the solution, the $P_2$ methods seems to be a good compromise between accuracy and efficiency. We notice that the error $err_3$ on amplitude between our solutions and the reference seems to converge towards a value equal to 4-5 %. It would be interesting to compare these results with those corresponding to a source introduced by a Dirac function, perhaps more suitable with high-order methods.

The results obtained, at station C1, with the method $P_2$ and a mesh spacing $h=200m$ are presented in figure 11. They are also compared to the reference solution. This constitutes a validation of the source implementation and of the absorbing boundary conditions since no spurious reflections from the boundaries appear in the profiles.

In a second step, we realize some additional simulations in order to compare the two types of absorbing boundary conditions ($\alpha = 1$ and 2), already studied in the previous test case. For this study, we choose the $P_4$ method and set the mesh size to $h = 300 \text{ m}$. Solutions are calculated until time 10.0 s using two computational domains: the initial one, which is $36\text{km} \times 36\text{km} \times 12\text{km}$, and a reduced one, $24\text{km} \times 24\text{km} \times 12\text{km}$. For this second domain, the boundary is 6km closer to the stations than for the initial one. We compare in Figure 12 the solutions for velocities $V_x$ and $V_z$ at station C1, on both domains and using both types of boundary conditions. When we examine the solutions for $V_x$ (figures a and b), we notice that reducing the domain produces some reflections, especially for times greater than 6.0s. The effect of the modified absorbing condition is noticeable on both figures, but only in the time interval [3; 5.5]. Comparing now the results for $V_z$, figures c and d, we find no difference due to the boundary conditions. Reflections are clearly visible and only explained by the distance to the boundary. If the modified absorbing condition allowed a clear improvement in the previous case, for a wave propagating perpendicular to the boundary, it does not reduce the reflections in this case, in presence of grazing waves, especially for the reduced domain.

6 Conclusion

We proposed a discontinuous Galerkin finite element method to solve the first-order hyperbolic system of the elastodynamic equations, written in velocity-stress formulation. This method is applied to tetraedral meshes allowing an accurate approximation of the medium or mesh refinement in particular areas of the domain. The
spatial approximation is based on high-degree polynomials defined locally on each element of the mesh. The scheme combines centered fluxes and a leap-frog scheme in time which leads to a non diffusive method. We also detailed a simple absorbing boundary condition, derived from an upwind scheme. In this paper, the properties of the method have been analysed in detail. First, we have proved that this explicit scheme is stable under a condition on the time-step resulting from a discrete energy which is preserved in domains with free surfaces and decreasing in presence of absorbing boundary conditions. Moreover, the convergence of the method for the semi-discrete and the fully discrete schemes has been studied. More precisely, for solutions regular enough, the scheme is convergent with an order $O(\Delta t^k + h^p)$ for domains with only free surface boundary conditions, or $O(\Delta t^{\min(p,2)} + h^k)$ when the absorbing boundary conditions are included in the scheme, where $k$ is the degree of the polynomial approximation, $N$ is the order of the leap-frog scheme and $\alpha \in \{1, 2, 3, 4\}$ is defined by (37).

Three-dimensional numerical results have been obtained with a parallel implementation of the DG method. The first one refers to the propagation of an eigenmode in a cubic cavity for which an exact solution is known and a numerical study of the convergence has been realised for series of uniform meshes, comparing the standard and the fourth-order leap-frog time scheme. In a second application, we study the propagation of a pulse which allows comparing two types of boundary conditions. Finally, the last test-case concerns the propagation of the waves due to an explosive source in a half-space which enables a validation of the source implementation and the absorbing boundary conditions. The main conclusions of this study are the accuracy and the efficiency of the method particularly for higher-degree polynomial interpolations. These results highlight the gain obtained with these methods, especially when compared to the standard finite volumes method ($P_0$) for which mesh refinement did not enable improving both errors on amplitudes and arrival time. The good compromise between accuracy and CPU time seems obtained using the P2 method (based on polynomials of degree 2). Moreover, the free surface condition is accurately approximated. These results are encouraging and could be developed in several interesting directions. First, the treatment at absorbing boundaries could be improved using more accurate boundary conditions. Many techniques have been proposed to reduce spurious reflections, including additional absorbing layers which may be costly in 3D. Then, high-order non reflecting boundary conditions [34] or an implicit treatment of our simple absorbing condition could be interesting ways of improvement, especially for grazing waves.

Finally, we are interested in simulations of more realistic three dimensional problems including topography and
realistic media.

Acknowledgement

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References


