

INVERSE PROBLEM FOR THE HELMHOLTZ EQUATION WITH CAUCHY DATA: RECONSTRUCTION WITH CONDITIONAL WELL-POSEDNESS DRIVEN ITERATIVE REGULARIZATION

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Abstract. In this paper, we study the performance of Full Waveform Inversion (FWI) from time-harmonic Cauchy data *via* conditional well-posedness driven iterative regularization. The Cauchy data can be obtained with dual sensors measuring the pressure and the normal velocity. We define a novel misfit functional which, adapted to the Cauchy data, allows the independent location of experimental and computational sources. The conditional well-posedness is obtained for a hierarchy of subspaces in which the inverse problem with partial data is Lipschitz stable. Here, these subspaces yield piecewise linear representations of the wave speed on given domain partitions. Domain partitions can be adaptively obtained through segmentation of the gradient. The domain partitions can be taken as a coarsening of an unstructured tetrahedral mesh associated with a finite element discretization of the Helmholtz equation. We illustrate the effectiveness of the iterative regularization through computational experiments with data in dimension three. In comparison with earlier work, the Cauchy data do not suffer from eigenfrequencies in the configurations.

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1. INTRODUCTION

Iterative methods for the recovery of subsurface parameters have been collectively referred to, in reflection seismology, as Full Waveform Inversion (FWI). FWI consists in the minimization of the residuals, defined as some difference between the observed and modelled data. FWI was originally formulated with time-domain data using an energy norm in the misfit functional by Lailly [23] and Tarantola [38, 39]. The time-harmonic formulation of the seismic inverse problem was later considered by Pratt *et al.* [32, 33]. Standardly, one applies the adjoint state method for the implementation of the minimization procedure.

Keywords and phrases. Inverse problems, Helmholtz equation, stability and convergence of numerical methods, reconstruction algorithm.

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To mitigate the nonlinearity and ill-posedness of the inverse problem, hierarchical multiscale strategies have been developed. In the time-domain, Bunks *et al.* [11] proposed successive inversion of data subsets of increasing frequency contents. This multiscale approach can be related to the subspace search method introduced in Kennett *et al.* [21]. Multiscale Gauss–Newton–Krylov methods were developed by Akcelik *et al.* [2] and many further developments have taken place since then. The application of wavelet bases enabling successive levels of model compression were considered by Loris *et al.* [26, 27] in wave-equation tomography and in FWI by Lin *et al.* in [25], Yuan and Simons in [43, 44]. These ideas were natural but lacked foundation and understanding of convergence. In De Hoop *et al.* [18], the authors developed an iterative regularization approach with a multilevel strategy derived from conditional Lipschitz stability estimates with a convergence analysis. For the regularization they introduced a multiscale hierarchy of subspaces for which the Lipschitz stability estimates hold, with an associated model compression rate through projections. This compression rate appears in the convergence analysis and mitigates the growth of stability constants with scale refinement, that is, increasing dimension of the subspaces. In the time-harmonic case, such quantitative stability estimates have been proved for piecewise constant [7] and piecewise linear [3] representations of the wave speed with a given domain partition; we refer to [3] for more extended bibliography. Note that, in view of recent work by Cârstea *et al.* [14], on a germane problem in elasticity, it seems quite possible to recover the partition as well, in the case where the domain is composed of sub-analytic sets. In the piecewise constant case, this includes subspaces defined by Haar wavelets, where the analysis in [3] is adapted to general domain partitions such as unstructured tetrahedral meshes, that can be associated with a segmentation. Here, we present a computational framework of our reconstruction algorithm *via* iterative regularization using piecewise linear representations as stable subspaces, as well as experiments.

We use Cauchy data assimilated from dual sensor acquisition. Cauchy data do not suffer from eigenfrequencies unlike the Dirichlet-to-Neumann map, which, in fact, cannot be observed directly in seismic marine acquisition. For a perspective on dual sensor acquisition devices and simultaneous measurements of pressure and vertical or normal velocity, we refer to [13, 40]. Dual sensor acquisition has additional benefits such as in noise reduction [34, 41]. Independent of earlier analysis, the results presented here confirm the relevance of acquiring dual sensor data. The marine seismic acquisition we consider consists in sources positioned above the fixed receivers lattice. This setup matches the TopSeis acquisition system that has been recently deployed by CGG (Compagnie Générale de Géophysique) and Lundin Norway AS⁵, for the purpose of improving the near offset data.

The formulation of FWI, and its underlying minimization procedure, has been proposed with norms different from energy norms. Shin and collaborators [6, 36, 37] compared the use of the phase and/or amplitude information in the data. Envelope based misfit functional is considered in [8, 42]. Brossier, Operto and Virieux investigate the use of the L^1 norm in [9], and in combination with L^2 norm in [10]. In this paper, we introduce a new misfit functional which is based upon conditional stability of the inverse problem for Cauchy data. This functional, related to the Green’s identity, is formulated in terms of repeated integration of a quadratic expression. Such a formulation overcomes the difficulties of computational complexity occurring in discretizing operator norms, or distances of subspaces (which are typically used in theoretical stability estimates, [3, 4]). Further, it allows independent locations for field and computational sources in the discrete settings.

The paper is organized as follows. We detail the inverse problem and the main stability results in Section 2, where the misfit functional to minimize for the reconstruction procedure is given. In Section 2.4 the computation of the gradient of the misfit functional is conducted using the Lagrangian approach. In Section 3, numerical experiments are presented to demonstrate the efficiency of the algorithm, using single-frequency Cauchy data. Finally, Section 4 shows an experiment with different locations of observational and computational sources.

⁵See <https://www.cgg.com/en/What-We-Do/Offshore/Products-and-Solutions/TopSeis>

2. MISFIT FUNCTIONAL AND STABILITY RESULT

2.1. Assumptions about the domain Ω

We fix some notations that will be adopted in this paper. Given a point $x \in \mathbb{R}^3$, with $x = (x', x_3)$, where $x' \in \mathbb{R}^2$ and $x_3 \in \mathbb{R}$, $B_r(x), B'_r(x')$ denote the open balls in $\mathbb{R}^3, \mathbb{R}^2$ centred at x, x' respectively with radius r . We also denote by $Q_r(x)$ the cylinder

$$Q_r(x) = B'_r(x') \times (x_3 - r, x_3 + r)$$

and $B_r = B_r(0)$ and $Q_r = Q_r(0)$.

- (1) We assume that Ω is a domain in \mathbb{R}^3 and that there exist positive constants r_0 and B (r_0 being dimensionally a length, B being an absolute constant), such that

$$|\Omega| \leq Br_0^3, \tag{2.1}$$

where $|\Omega|$ denotes the Lebesgue measure of Ω .

- (2) We fix an open non-empty subset Σ of $\partial\Omega$ (where the measurements in terms of the local Cauchy data are taken).
- (3) We assume that Ω can be decomposed as follows:

$$\bar{\Omega} = \bigcup_{j=1}^N \bar{D}_j,$$

where $D_j, j = 1, \dots, N$ are known open sets of \mathbb{R}^3 , satisfying the conditions below.

- (a) $D_j, j = 1, \dots, N$ are connected and pairwise non-overlapping polyhedrons.
- (b) $\partial D_j, j = 1, \dots, N$ are of Lipschitz class with constants r_0, L see, for example, [1]).
- (c) There exists one region, say D_1 , such that $\partial D_1 \cap \Sigma$ contains a *flat* portion Σ_1 of size r_0 and for every $j \in \{2, \dots, N\}$, we can select a subchain $\{D_{j_k}\}_{k=1}^K$ ($K \leq N$) of the partition $\{D_j\}_{j=1}^N$ of Ω , such that

$$D_{j_1} = D_1, \quad D_{j_K} = D_j. \tag{2.2}$$

In addition, for every $k = 1, \dots, K, D_{j_{k-1}}$ and D_{j_k} are contiguous in the sense that $\partial D_{j_k} \cap \partial D_{j_{k-1}}$ contains a *flat* portion Σ_k of size r_0 (here we agree that $D_{j_0} = \mathbb{R}^3 \setminus \Omega$), such that

$$\Sigma_k \subset \Omega, \quad \text{for every } k = 2, \dots, K.$$

We illustrate the configuration and selection of subchain in Figure 1. We emphasize that under such an assumption, for every $k = 1, \dots, K$, there exists $P_k \in \Sigma_k$ and a rigid transformation of coordinates (depending on k) under which we have $P_k = 0$ and

$$\begin{aligned} \Sigma_k \cap Q_{r_0/3} &= \{x \in Q_{r_0/3} \mid x_3 = 0\}, \\ D_{j_k} \cap Q_{r_0/3} &= \{x \in Q_{r_0/3} \mid x_3 > 0\}, \\ D_{j_{k-1}} \cap Q_{r_0/3} &= \{x \in Q_{r_0/3} \mid x_3 < 0\}. \end{aligned} \tag{2.3}$$

2.2. *A priori* information on the wave speed c

We shall consider a real valued function $c \in L^\infty(\Omega)$, with

$$0 < c_{\min} \leq c \leq c_{\max}, \tag{2.4}$$

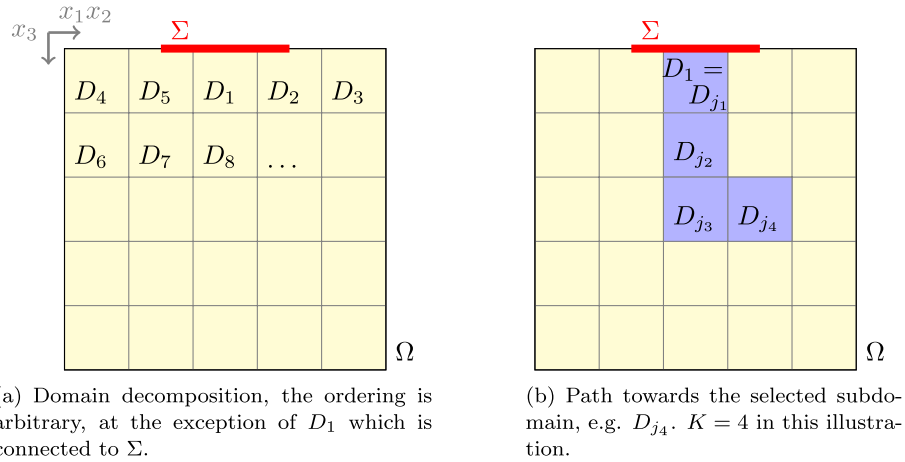


FIGURE 1. Illustration of the domain decomposition and the path towards a selected subdomain D_{jk} .

for some positive constants c_{\min} , c_{\max} and of type

$$c(x) = \sum_{j=1}^N c_j(x)\chi_{D_j}(x), \quad x \in \Omega, \tag{2.5a}$$

$$c_j(x) = a_j + A_j \cdot x, \tag{2.5b}$$

where $a_j \in \mathbb{R}$, $A_j \in \mathbb{R}^3$ are scalar and vector valued constants respectively, and $\{D_j\}_{j=1}^N$, $j = 1, \dots, N$ is the given partition of Ω introduced in Section 2.1.

2.3. Cauchy data and misfit functional

We let Ω denote the subsurface domain with Lipschitz boundary $\partial\Omega$ and Σ the non-empty open portion of $\partial\Omega$ introduced in Section 2.1 where the acquisition is carried out. We introduce the subspace of $H^{1/2}(\partial\Omega)$,

$$H_{co}^{1/2}(\Sigma) = \{f \in H^{1/2}(\partial\Omega) \mid \text{supp } f \subset \Sigma\}.$$

Its closure with respect to the $H^{1/2}(\partial\Omega)$ norm is the space $H_{00}^{1/2}(\Sigma)$. In a similar manner, we define $H_{00}^{1/2}(\partial\Omega \setminus \overline{\Sigma})$. We denote the pressure by u . The *Cauchy data* associated to c is the space \mathcal{C}_c^Σ ,

$$\mathcal{C}_c^\Sigma = \left\{ (f, g) \in H^{1/2}(\partial\Omega)|_\Sigma \times H^{-1/2}(\partial\Omega)|_\Sigma \mid \exists u \in H^1(\Omega) \text{ weak solution to } \Delta u + k^2 c^{-2} u = 0 \text{ in } \Omega, \right. \\ \left. u|_\Sigma = f, \right. \\ \left. \left\langle \frac{\partial u}{\partial \nu} \Big|_{\partial\Omega}, \varphi \right\rangle = \langle g, \varphi \rangle, \forall \varphi \in H_{00}^{1/2}(\Sigma) \right\}.$$

Here, $\langle \psi, \varphi \rangle$ denotes the duality between the complex valued spaces $H^{-\frac{1}{2}}(\partial\Omega)$, $H^{\frac{1}{2}}(\partial\Omega)$ based on the L^2 inner product

$$\langle \psi, \varphi \rangle = \int_{\partial\Omega} \psi \overline{\varphi}$$

and $H^{\frac{1}{2}}(\partial\Omega)|_{\Sigma}$ and $H^{-\frac{1}{2}}(\partial\Omega)|_{\Sigma}$ denote the *restrictions* of $H^{\frac{1}{2}}(\partial\Omega)$ and $H^{-\frac{1}{2}}(\partial\Omega)$ to Σ respectively. \mathcal{C}_c^{Σ} is a subspace of the Hilbert space $H^{\frac{1}{2}}(\partial\Omega)|_{\Sigma} \times H^{-\frac{1}{2}}(\partial\Omega)|_{\Sigma}$.

We embed Ω in an ambient domain $\mathcal{Y} \supset \Omega$ as we will find convenient to introduce Green’s function not precisely for the physical domain Ω but for this augmented domain \mathcal{Y} .

We recall that by assumption 3(c) of Section 2.1 we can assume that there exists a point P_1 such that up to a rigid transformation of coordinates we have that $P_1 = 0$ and (2.3) holds with $\Sigma = \Sigma_1$. Denoting by

$$D_0 = \left\{ x \in (\mathbb{R}^3 \setminus \Omega) \cap B_{\frac{r_0}{3}} \mid |x_i| < \frac{r_0}{6}, i = 1, 2; -\frac{r_0}{6} < x_3 < 0 \right\},$$

it turns out that the augmented domain $\mathcal{Y} = \overline{(\Omega \cup D_0)}^{\circ}$ is of Lipschitz class with constants $\frac{r_0}{3}$ and \tilde{L} , where \tilde{L} depends on L only. Given $r > 0$, we set

$$\Gamma_1 = \left\{ x \in \mathcal{Y} \mid |x_i| < \frac{r_0}{6}, i = 1, 2; x_3 = -\frac{r_0}{6} \right\}, \tag{2.6}$$

$$\Gamma_2 = \partial\mathcal{Y} \setminus \bar{\Gamma}_1, \tag{2.7}$$

$$(\mathcal{Y})_r = \left\{ x \in \mathcal{Y} \mid \text{dist}(x, \partial\mathcal{Y}) > r \right\}. \tag{2.8}$$

We also introduce the following sets

$$D'_0 = \left\{ x \in D_0 \mid -\frac{r_0}{12} < x_3 < 0 \right\}, \tag{2.9}$$

$$D''_0 = \left\{ x \in D_0 \mid -\frac{r_0}{6} < x_3 < -\frac{r_0}{12} \right\}, \tag{2.10}$$

$$K_0 = \left\{ x \in D''_0 \mid |x_i| < \frac{r_0}{24}, i = 1, 2; -\frac{13}{96}r_0 < x_3 < -\frac{11}{96}r_0 \right\}, \tag{2.11}$$

$$K_1 = \left\{ x \in D''_0 \mid |x_i| < \frac{r_0}{12}, i = 1, 2; -\frac{7}{48}r_0 < x_3 < -\frac{5}{48}r_0 \right\}, \tag{2.12}$$

$$\Omega' = \overline{(D'_0 \cup \Omega)}^{\circ}. \tag{2.13}$$

Note that, fixing the origin at the center of D''_0 , K_0, K_1 are concentric parallelograms scaled by the factors $\frac{1}{4}, \frac{1}{2}$ respectively. We illustrate the geometry in Figure 2. Note also that this precise choice of scale parameters is just made for the sake of definiteness. What really matters is the general geometrical configuration, in particular we must have

$$\text{dist}(D''_0, \bar{\Omega}) \geq \frac{r_0}{12} > 0 \quad \text{and} \quad K_0 \subset\subset K_1 \subset\subset D''_0.$$

We shall denote by $\Gamma(x, y)$ the standard fundamental solution to the Laplace equation which is

$$\Gamma(x, y) = \frac{1}{4\pi|x - y|}. \tag{2.14}$$

Definition 2.1. Let B, N, r_0, L, c_{\min} and c_{\max} be given positive numbers with $N \in \mathbb{N}$. We will refer to this set of numbers as to the *a priori data*. Several constants depending on the *a priori data* will appear within the paper. In order to simplify our notation, any quantity denoted by C, C_1, C_2, \dots will be called a *constant* understanding in most cases that it only depends on the *a priori data*.

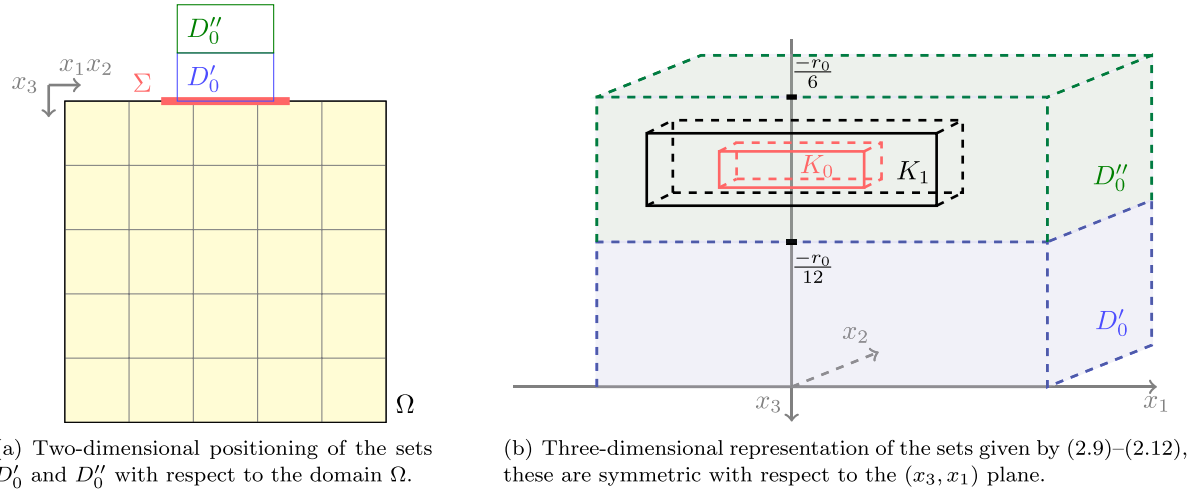


FIGURE 2. Illustration of the sets that are defined in (2.9)–(2.12).

Next we introduce a mixed boundary value problem for $\Delta + k^2 c^{-2}$ in \mathcal{Y} which is always well posed, independently of any *a priori* condition on c , besides the assumption of being real valued and bounded. This shall enable us to construct Green’s function for $\Delta + k^2 c^{-2}$ in \mathcal{Y} .

We assume that the uppermost part D_0 of the domain \mathcal{Y} represents a region filled by water. The wave speed in D_0 can then be assumed to be known and constant with

$$c = c_0, \quad \text{in } D_0. \tag{2.15}$$

Pressure sources (air gun) are excited to produce impulses located at points in $K_1 \subset D_0$ and Cauchy data are collected through dual sensors located at the surface Σ , which lies below K_1 .

Pressure is assumed to be zero at the sea level Γ_1 (*i.e.*, free surface), and to satisfy, on the remaining part of the boundary Γ_2 (of the region of interest), a (conventional) absorbing condition. If we model this problem in the frequency domain, and assume that the source is modeled by a Dirac’s delta concentrated at a point y , the pressure is represented as the Green’s function of the following mixed boundary value problem,

$$\begin{cases} \Delta G_c(\cdot, y) + k^2 c^{-2}(\cdot) G_c(\cdot, y) = -\delta(\cdot - y), & \text{in } \mathcal{Y}, \\ G_c(\cdot, y) = 0, & \text{on } \Gamma_1, \\ \partial_\nu G_c(\cdot, y) - ik_0 G_c(\cdot, y) = 0, & \text{on } \Gamma_2. \end{cases} \tag{2.16}$$

The theory developed, for example, in [12] shows that such a function G_c exists and is unique in the case of constant wave speed c . Note that the term k_0 is conventionally assumed to be constant and known. The next theorem collects the main features of the Green’s function solving (2.16) also in the case of variable wave speed $c(x)$. A similar result, but with stronger hypothesis, was proven in [3, Props. 3.1, 3.4, 3.5]. The thesis here is slightly weaker, but the argument is somewhat simpler.

Theorem 2.2. *For any $y \in \mathcal{Y}$, there exists a unique distributional solution $G_c(\cdot, y)$ to (2.16). Moreover, there exists a constant $C > 0$ depending on r_0, L, k and on c_{\min} such that for any $x, y \in (\mathcal{Y})_{r_0}$, $x \neq y$ we have that*

$$|G_c(x, y) - \Gamma(x, y)| \leq C \tag{2.17}$$

and

$$|\nabla_y G_c(x, y) - \nabla_y \Gamma(x, y)| \leq C(|\log |x - y|| + 1). \tag{2.18}$$

Moreover, let Q_{k+1} be a point such that $Q_{k+1} \in B_{\frac{r_0}{8}}(P_{k+1}) \cap \Sigma_{k+1}$ with $k \in \{1, \dots, N - 1\}$, then the following inequality holds true for every $x \in B_{\frac{r_0}{16}}(P_{k+1}) \cap D_{j_{k+1}}$ and every $y = Q_{k+1} - re_3$, where $r \in (0, \frac{r_0}{16})$

$$|\nabla_y^2(G_c(x, y) - \Gamma(x, y))| \leq Cr^{-1}. \tag{2.19}$$

Here ∇_y^2 denotes the Hessian matrix.

Proof. Let $y \in \mathcal{Y}$ and let G_0 be the Green’s function for the Laplace operator which solves

$$\begin{cases} \Delta G_0(\cdot, y) = -\delta(\cdot - y), & \text{in } \mathcal{Y}, \\ G_0(\cdot, y) = 0, & \text{on } \Gamma_1, \\ \partial_\nu G_0(\cdot, y) - ik_0 G_0(\cdot, y) = 0, & \text{on } \Gamma_2. \end{cases} \tag{2.20}$$

The existence and uniqueness of a distributional solution $G_0 \in L^1(\mathcal{Y})$ to (2.20) is a consequence of standard theory on boundary value problems for the Laplace equation. By standard techniques it can be proved that for any $y \in \mathcal{Y}$ such that $\text{dist}(y, \partial\mathcal{Y}) \geq \frac{r_0}{4}$ we have that $G_0(\cdot, y) \in L^2(\mathcal{Y})$. Now we define $R_c(\cdot, y) \in H^1(\mathcal{Y})$ to be the solution to

$$\begin{cases} \Delta R_c(\cdot, y) + k^2 c^{-2} R_c(\cdot, y) = -k^2 c^{-2} G_0(\cdot, y), & \text{in } \mathcal{Y}, \\ R_c(\cdot, y) = 0, & \text{on } \Gamma_1, \\ \partial_\nu R_c(\cdot, y) - ik_0 R_c(\cdot, y) = 0, & \text{on } \Gamma_2. \end{cases} \tag{2.21}$$

The existence and uniqueness for (2.21) follows along the lines of the proof of [3], Proposition 3.1, which relies on the Fredholm alternative theory. Moreover, by arguments based on well-known estimates for the Cauchy problem contained in Proposition 3.1 from [3], we have that

$$\|R_c(\cdot, y)\|_{H^1(\mathcal{Y})} \leq C \|G_0(\cdot, y)\|_{L^2(\mathcal{Y})} \tag{2.22}$$

and, by standard interior estimates, that

$$|R_c(x, y)| \leq C, \tag{2.23}$$

for any $x \neq y$, $x \in \mathcal{Y}$ and $\text{dist}(y, \partial\mathcal{Y}) \geq \frac{r_0}{4}$. If we form

$$G_c(x, y) = G_0(x, y) + R_c(x, y), \tag{2.24}$$

then we end up with the following estimate

$$|G_c(x, y)| \leq C|x - y|^{-1}, \tag{2.25}$$

for any $x, y \in \mathcal{Y}$, $x \neq y$ and $\text{dist}(x, \partial\mathcal{Y}) \geq \frac{r_0}{4}$, $\text{dist}(y, \partial\mathcal{Y}) \geq \frac{r_0}{4}$. The latter combined with the arguments in the proofs of Propositions 3.4, 3.5 from [3] suffices to deduce (2.17), (2.18) and (2.19). \square

Assuming that, for sources placed at arbitrary points $z \in K_1$, we can measure associated Cauchy data on Σ :

$$G_{\text{obs}}(x, z), \quad \frac{\partial}{\partial \nu_x} G_{\text{obs}}(x, z), \quad x \in \Sigma, \tag{2.26}$$

we seek c which minimizes the following misfit functional

$$\mathcal{J}(c) = \int_{K_1 \times K_1} \left| \int_{\Sigma} (G_c(x, y) \partial_\nu G_{\text{obs}}(x, z) - G_{\text{obs}}(x, z) \partial_\nu G_c(x, y)) \, d\mu(x) \right|^2 \, dy \, dz, \tag{2.27}$$

where μ denotes the element of surface measure.

The introduction of the misfit functional (2.27) is motivated by the following argument. Given two wave speeds $c^{(i)}$, $i = 1, 2$, consider the Green's functions G_i , introduced in Theorem 2.2, corresponding to $c^{(i)}$ in \mathcal{Y} and the following quantity

$$\mathcal{J}(c^{(1)}, c^{(2)}) = \int_{K_1 \times K_1} |\mathcal{S}u_0(y, z)|^2 dy dz, \quad (2.28)$$

where

$$\mathcal{S}u_0(y, z) = \int_{\Sigma_1} (G_1(x, y) \partial_\nu G_2(x, z) - G_2(x, z) \partial_\nu G_1(x, y)) d\mu(x), \quad \text{for any } y, z \in D_0. \quad (2.29)$$

Expressions of the form above have appeared in many occasions in the treatment of inverse boundary problems. Analogies can be found with the probe method by Ikehata [20], see also [31]. In particular, a very strong relation can be observed with the so-called *reciprocity gap functional* introduced by Colton and Haddar [17] for inverse scattering.

It would be a matter of an exercise to show that, fixing $y \in K_1$,

$$c \in L^\infty(\mathcal{Y}) \rightarrow (G_c(x, y)|_\Sigma, \partial_\nu G_c(x, y)|_\Sigma) \in H^{1/2}(\partial\Omega)|_\Sigma \times H^{-1/2}(\partial\Omega)|_\Sigma \quad (2.30)$$

is Fréchet differentiable. Note also that, since we are assuming (2.5a) and (2.5b) (that is, c lives in a finite dimensional space), the $L^\infty(\mathcal{Y})$ norm can equivalently be replaced by the $L^2(\mathcal{Y})$ norm. This will enable us to apply to \mathcal{J} a projected steepest descent method in Section 3.

Theorem 2.3. *Let Ω , D_j , $j = 1, \dots, N$ and Σ be a domain, N subdomains of Ω and a portion of $\partial\Omega$ as in section 2.1 respectively. Let $c^{(i)}$, $i = 1, 2$ be two wave speeds satisfying (2.4) and of type*

$$c^{(i)} = \sum_{j=1}^N c_j^{(i)}(x) \chi_{D_j}(x), \quad x \in \Omega, \quad (2.31)$$

where

$$c_j^{(i)}(x) = a_j^{(i)} + A_j^{(i)} \cdot x,$$

with $a_j^{(i)} \in \mathbb{R}$ and $A_j^{(i)} \in \mathbb{R}^3$, then we have

$$\|c^{(1)} - c^{(2)}\|_{L^\infty(\Omega)} \leq C \left(\mathcal{J}(c^{(1)}, c^{(2)}) \right)^{1/2}, \quad (2.32)$$

where C is a positive constant that depends on the a priori data only.

Remark 2.4. Note that the introduction of the misfit functional is driven on the one hand by our computational experiments and on the other hand, it is inspired by the method of singular solutions used in previous stability results (see [3] for the case of the Helmholtz equation). Although a natural metric on the space of Cauchy data \mathcal{C}_c^Σ is given by the *distance (aperture) d* introduced in equation (2.8) from [3], we have

$$\mathcal{J}(c^{(1)}, c^{(2)}) \leq Cd(\mathcal{C}_{c^{(1)}}^\Sigma, \mathcal{C}_{c^{(2)}}^\Sigma)^2 \quad (2.33)$$

and also

$$\mathcal{J}(c^{(1)}, c^{(2)}) \leq C \|A_{c^{(1)}}^\Sigma - A_{c^{(2)}}^\Sigma\|_*^2, \quad (2.34)$$

where $A_c^\Sigma : u \in H^{1/2}(\partial\Omega)|_\Sigma \rightarrow \partial_\nu u \in H^{-1/2}(\partial\Omega)|_\Sigma$ is the local Dirichlet-to-Neumann map with its natural norm (here denoted by $\|\cdot\|_*$) between local trace spaces. These estimates justify the use of \mathcal{J} as a substitute to the more traditional quantifications of the error on boundary data (either involving distance of spaces of Cauchy data or boundary maps).

Proof of Theorem 2.3.

The proof requires only some adaptations of Theorem 2.2 from [3], which are outlined below.

- i) We introduce different boundary conditions. This aspect involves some modifications in the constructions of the Green’s function and has been treated in Theorem 2.2 above. Note that we took advantage of the fact that here we focus on the three-dimensional case only.
- ii) We replace the domain Ω with Ω' and add as initial subdomain D'_0 instead of D_1 . We take advantage of the fact that $c^{(1)} = c^{(2)} = 1$ in $D'_0 \subset \mathcal{Y} \setminus \bar{\Omega}$ which allows us to skip the arguments of [3, Sect. 4.3].
- iii) We observe that in Theorem 2.2 from [3] the right-hand side in formula (2.20) is expressed in terms of the distance between spaces of Cauchy data. However, the only Cauchy data that are actually used are those arising from Green’s function with pole in \widehat{D}_0 . And the role of \widehat{D}_0 (see [3]) can be equivalently taken by the set K_0 introduced here in (2.11). Moreover, such Cauchy data intervene only in expressions like the one in (2.29) above. Therefore, the right-hand side of [3, Eq. (2.20)] can be replaced by

$$\sup_{K_0 \times K_0} |S_{\mathcal{U}_0}(y, z)|. \tag{2.35}$$

We also recall that $S_{\mathcal{U}_0}(y, z)$ is a solution to

$$(\Delta_y + \Delta_z + 2k^2)S_{\mathcal{U}_0}(y, z) = 0 \tag{2.36}$$

in $D_0 \times D_0$ (see [3], Eq. (4.61)). Consequently, by standard estimate of boundedness in the interior we have

$$\sup_{K_0 \times K_0} |S_{\mathcal{U}_0}(y, z)|^2 \leq C \int_{K_1 \times K_1} |S_{\mathcal{U}_0}(y, z)|^2 \, dy \, dz \tag{2.37}$$

where $C > 0$ is a constant depending on k and on r_0 .

- iv) Another difference comes from the fact that we are now assuming $c^{(i)}$ piecewise linear instead of $q^{(i)} = k^2(c^{(i)})^{-2}$, for $i = 1, 2$. However due to the assumption (2.4) the estimation of $c^{(i)}, \nabla c^{(i)}$ at each interface is equivalent to that for $q^{(i)}, \nabla q^{(i)}$, for $i = 1, 2$.

□

The stability given in Theorem 2.3 justifies the use of (2.27) for an optimization algorithm toward the inversion of *Cauchy data to wave speed*.

2.4. Computation of the gradient using the adjoint state method

We start by observing that, although $c \rightarrow G_c(x, y)$ does not map into $H^1(\mathcal{Y})$ (because of the singularity of the Green’s function), the derivative $D_c G_c(x, y)[c]\delta c$ exists and does belong in $H^1(\mathcal{Y})$. This can be achieved by recalling (2.21) and (2.24). That is

$$R_c(x, y) = G_c(x, y) - G_0(x, y) \in H^1(\mathcal{Y}), \tag{2.38}$$

and the second term G_0 is independent of c . Hence we may well define

$$D_c G_c(x, y)[c]\delta c = D_c R_c(x, y)[c]\delta c. \tag{2.39}$$

We denote

$$H^1_{\Gamma_1}(\mathcal{Y}) = \{v \in H^1(\mathcal{Y}) \mid v|_{\Gamma_1} = 0\}. \tag{2.40}$$

We continue with the variational formulation of problem (2.16). As noted above, denoting $R_c(x, y) = G_c(x, y) - G_0(x, y)$, it can be formulated as

$$\int_{\mathcal{Y}} k^2 c^{-2}(x)G_c(x, y)v(x) - \nabla_x R_c(x, y) \cdot \nabla_x v(x) \, dx + \int_{\Gamma_2} ik_0 R_c(x, y)v(x) \, d\mu(x) = 0, \tag{2.41}$$

for every $v \in H^1_{\Gamma_1}(\mathcal{Y})$.

The parameter reconstruction is conducted *via* an iterative minimization of the misfit functional \mathcal{J} of (2.27), in a gradient descent algorithm. Therefore, we require the computation of the gradient of \mathcal{J} . For this purpose, we employ the adjoint state method, which allows the computation of the gradient without having to form explicitly the derivative G'_c . The method arose from the work of Lions and Mitter [24] and was promoted in the context of parameter derivation in Chavent [15]. It has massively been employed since then, and we refer to Plessix [30] for a review in the geophysical framework. Here, we follow the traditional steps for the selection of the Lagrange multiplier and formation of the gradient which are detailed, for example, in [16, 22], and that we adapt to our choice of misfit functional.

We first postpone the sum over the sources in the misfit functional (2.27), and select a single source for G_c and G_{obs} , y and z respectively. We introduce

$$\mathcal{I}(c)(y, z) = \left| \int_{\Sigma} \left(G_c(x, y) \partial_{\nu} G_{\text{obs}}(x, z) - G_{\text{obs}}(x, z) \partial_{\nu} G_c(x, y) \right) d\mu(x) \right|^2 = |S_{\mathcal{U}_0}(y, z)|^2, \tag{2.42}$$

such that

$$\mathcal{J}(c) = \int_{K_1 \times K_1} \mathcal{I}(c)(y, z) dy dz. \tag{2.43}$$

The Riesz representation theorem gives

$$D_c \mathcal{I}[c] \delta c = \int_{\mathcal{Y}} \nabla_c \mathcal{I} \delta c dx, \tag{2.44}$$

where $\nabla_c \mathcal{I}$ is the gradient, and $D_c \mathcal{I}$ is the differential defined for every $\delta c \in L^2(\mathcal{Y})$ by

$$D_c \mathcal{I}[c] : \delta c \rightarrow \lim_{h \rightarrow 0} \frac{\mathcal{I}(c + h\delta c) - \mathcal{I}(c)}{h}. \tag{2.45}$$

The adjoint state method considers the constrained minimization problem

$$\min_c \mathcal{I}(c) \quad \text{subject to (2.16)}. \tag{2.46}$$

The constraint can be replaced by the variational formulation (2.41) and the associated formulation of the Lagrangian is defined by

$$\begin{aligned} \mathcal{L}(c, G, \tilde{\gamma}) &= \mathcal{I}(c)(y, z) + \int_{\mathcal{Y}} k^2 c^{-2}(x) G(x, y) \tilde{\gamma}(x, y, z) - \nabla_x R(x, y) \cdot \nabla_x \tilde{\gamma}(x, y, z) dx \\ &+ \int_{\Gamma_2} ik_0 R(x, y) \tilde{\gamma}(x, y, z) d\mu(x). \end{aligned} \tag{2.47}$$

Here, $\tilde{\gamma}$ has the role of a Lagrange multiplier and a specific choice of it, γ , will be specified later. By letting $G = G_c$ be the solution of the forward problem, and hence $R = R_c$, we can form the Fréchet derivative. For the sake of brevity, we use the following notation

$$G' = D_c G_c[c] \delta c, \tag{2.48}$$

and we omit the variables x, y, z (keeping in mind that G_{obs} depends on (x, z) , G_c and G' on (x, y) , and $\tilde{\gamma}$ on (x, y, z)).

$$\begin{aligned} D_c \mathcal{I}[c] \delta c &= \text{Re} \left(D_c \mathcal{L}(c, G, \tilde{\gamma}) \delta c \right) \Big|_{G=G_c} \\ &= \text{Re} \left(2\overline{S_{\mathcal{U}_0}} \int_{\Sigma} (G' \partial_{\nu_x} G_{\text{obs}} - G_{\text{obs}} \partial_{\nu_x} G') d\mu(x) \right. \\ &+ \int_{\mathcal{Y}} \left(k^2 c^{-2} G' \tilde{\gamma} - 2k^2 c^{-3} G_c \tilde{\gamma} \delta c - \nabla_x G' \cdot \nabla_x \tilde{\gamma} \right) dx \\ &\left. + \int_{\Gamma_2} ik_0 G' \tilde{\gamma} d\mu(x) \right). \end{aligned} \tag{2.49}$$

Grouping together all the terms containing G' and replacing it by an arbitrary test function $v \in H_{\Gamma_1}^1(\mathcal{Y})$, the adjoint state γ is chosen as the solution to

$$2\overline{S_{U_0}} \int_{\Sigma} (v\partial_{\nu_x} G_{\text{obs}} - G_{\text{obs}}\partial_{\nu_x} v) \, d\mu(x) + \int_{\mathcal{Y}} (k^2 c^{-2} v \gamma - \nabla_x v \cdot \nabla_x \gamma) \, dx \tag{2.50}$$

$$+ \int_{\Gamma_2} ik_0 v \gamma \, d\mu(x) = 0. \tag{2.51}$$

Note that the first term,

$$v \rightarrow 2\overline{S_{U_0}} \int_{\Sigma} (v\partial_{\nu_x} G_{\text{obs}} - G_{\text{obs}}\partial_{\nu_x} v) \, d\mu(x), \tag{2.52}$$

is a bounded linear functional of $H_{\Gamma_1}^1(\mathcal{Y})$. Hence, by the arguments already mentioned in [3, Prop. 3.1], there exists a unique solution $\gamma \in H_{\Gamma_1}^1(\mathcal{Y})$ to problem (2.50).

With this choice of adjoint state, (2.49) reduces to

$$D_c \mathcal{I}[c] \delta c = \text{Re} \left(\int_{\mathcal{Y}} -2k^2 c^{-3} G_c \gamma \delta c \, dx \right). \tag{2.53}$$

Reassembling with $(y, z) \in K_1 \times K_1$, we get

$$\nabla_c \mathcal{J}(x) = -\text{Re} \left(\int_{K_1 \times K_1} 2k^2 c^{-3}(x) G_c(x, y) \gamma(x, y, z) \, dy \, dz \right). \tag{2.54}$$

Note that G_c is independent of z , hence, posing

$$\widehat{\gamma}(x, y) = \int_{K_1} \gamma(x, y, z) \, dz, \tag{2.55}$$

we have that $\widehat{\gamma}$ verifies, for every $v \in H_{\Gamma_1}^1(\mathcal{Y})$,

$$2 \int_{K_1} \overline{S_{U_0}} \int_{\Sigma} (v\partial_{\nu_x} G_{\text{obs}} - G_{\text{obs}}\partial_{\nu_x} v) \, d\mu(x) \, dz + \int_{\mathcal{Y}} (k^2 c^{-2} v \widehat{\gamma} - \nabla_x v \cdot \nabla_x \widehat{\gamma}) \, dx + \int_{\Gamma_2} ik_0 v \widehat{\gamma} \, d\mu(x) = 0. \tag{2.56}$$

3. COMPUTATIONAL EXPERIMENTS

For the computational experiments, the space coordinates will be denoted by (x, y, z) instead of (x_1, x_2, x_3) . Let us emphasize that the z coordinate is conventionally (in geophysical setup) oriented downwards and can be seen as the depth of the medium. We first consider a three-dimensional model (courtesy Statoil), which is illustrated in Figure 3. To have a clear visualization of the wave speed structures, we show horizontal and vertical sections at $z = 800$ m and $y = 1.125$ km respectively. The model is of size $2.55 \times 1.45 \times 1.22$ km with variations of wave speed from 1500 to 5200 m s^{-1} . We assume that the density is constant with $\rho = 1000$ kg m^{-3} .

The seismic acquisition consists of 160 sources and 1376 receivers, that is, dual sensors recording the Cauchy data probed by the sources. The receivers are positioned on a regular 43 (along the x -axis) by 32 (along the y -axis) grid at a fixed depth below the sources lattice. The configuration is illustrated in Figure 4. We consider two situations for the discretized sources map: the sources are first contained in a solid region, in accordance with the above analysis (see Fig. 4b). Then, they are restricted on a two-dimensional lattice (see Fig. 4c). The first approach, less common in seismic applications, is yet possible with recent acquisition technique described in Footnote 5, for which the depth of the sources can vary. Following the situation prescribed in Section 2.3, we assume that the uppermost part of the model (which is water), in which the Cauchy data are obtained, is

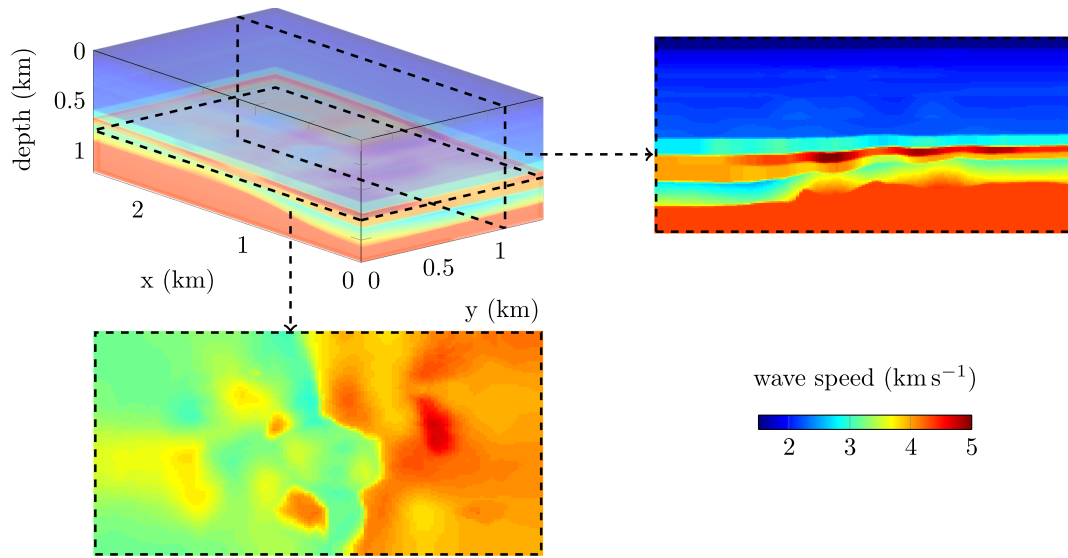


FIGURE 3. Three-dimensional representation, horizontal section at 800 m depth and vertical section at $y = 1.125$ km of the reference wave speed. It is represented with 1 527 168 nodal values.

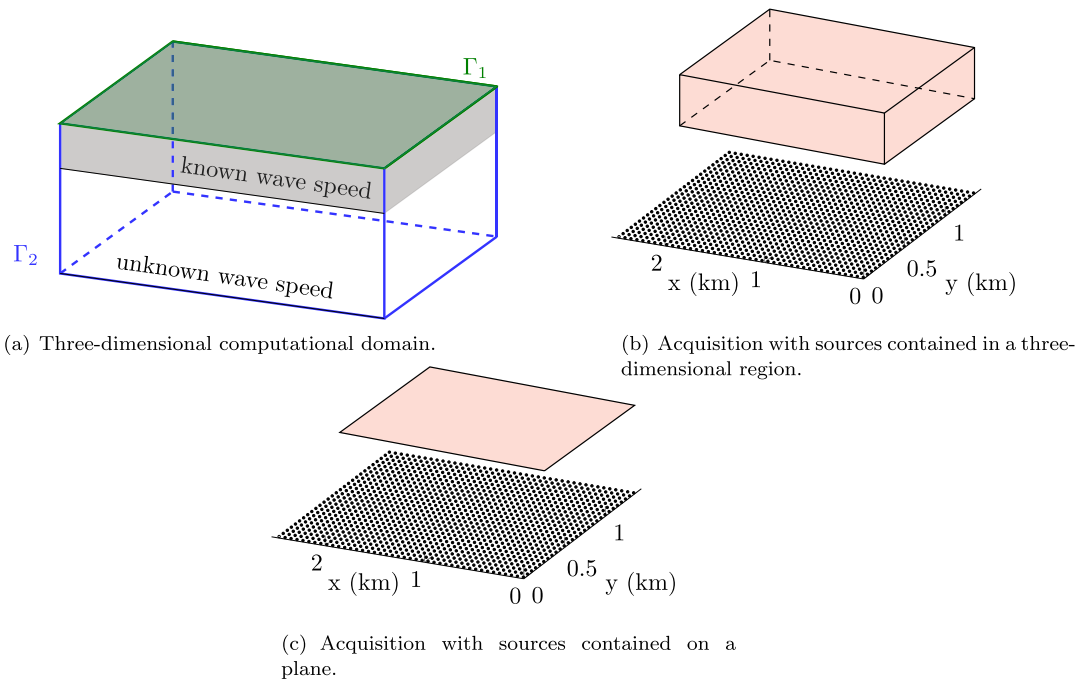


FIGURE 4. Illustration of the configuration. *Panel a:* we apply a Dirichlet boundary on the upper surface Γ_1 (in green) and absorbing boundary conditions on $\Gamma_2 = \partial\mathcal{Y} \setminus \bar{\Gamma}_1$ (indicated in blue). *Panels b–c:* sources that probe the Cauchy data lie in between the receivers and the free surface. Sources are positioned in a three or two-dimensional region (in red). The position of the receivers (black dots) remains fixed. Both receivers and sources lie in the area of known wave speed (Color online).

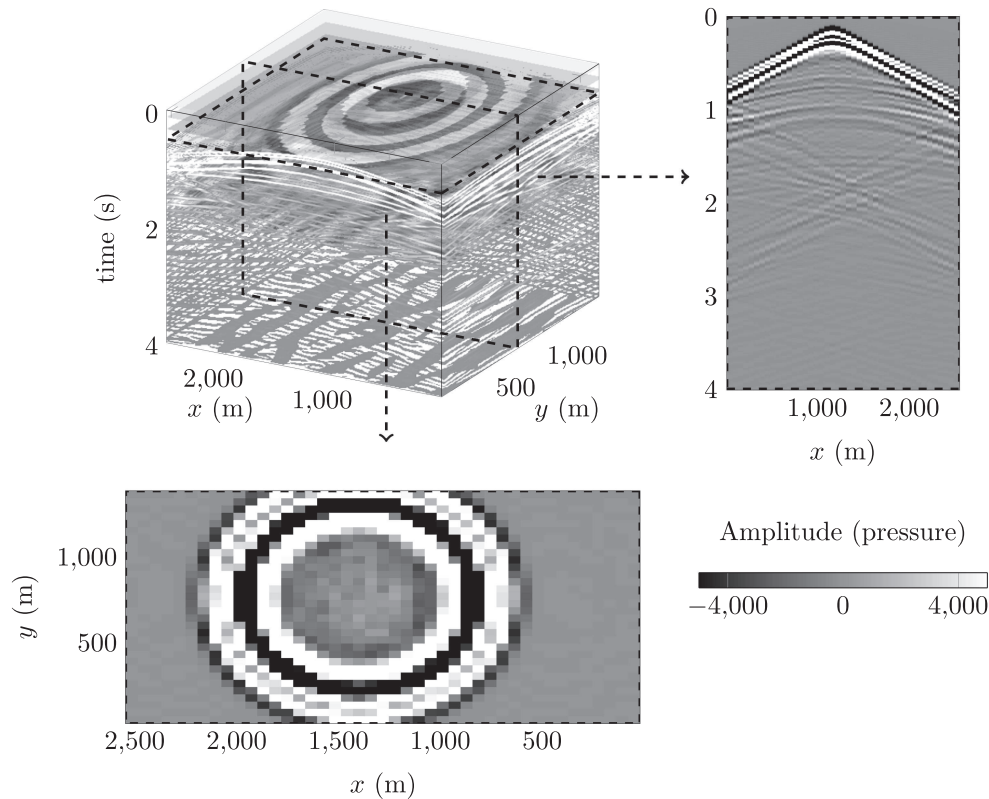


FIGURE 5. Three-dimensional time-domain pressure trace for a centrally located source. Sections at fixed time $t = 0.5\text{s}$ and for a fixed line of receivers positioned in $y = 695\text{m}$ are respectively given at the bottom and right of the 3D visualization. The x and y axis correspond with the receivers map (*i.e.*, every 60 m for x and 45 m for y).

known prior to the reconstruction. However, we do not assume the knowledge of the wave speed onto the lateral and bottom boundaries.

We impose a free surface, Dirichlet boundary condition on the top part, Γ_1 , of $\partial\mathcal{Y}$ and absorbing boundary conditions on $\Gamma_2 = \partial\mathcal{Y} \setminus \bar{\Gamma}_1$, by taking $k_0 = kc_0^{-1}$ in the third equation of (2.16), following Engquist and Majda [19].

Synthetic dual-sensor data are generated in the time-domain using a Discontinuous Galerkin (DG) finite element method⁶. The original data for a single centered source are presented in Figures 5 and 6 for the pressure and vertical velocity respectively. In the figures, we can observe the difference of scale in the amplitudes between the pressure and the vertical velocity.

We subject the (time-domain) data to Gaussian white noise, using a signal-to-noise ratio of 15 dB (this process is illustrated Fig. 7). Note that every receiver for each source has an independent white noise signal added. We apply the Fourier transform to these noisy data and obtain time-harmonic data. The effect of noise affects particularly the low-frequency regime in seismic, and frequencies below 3 Hz are usually unusable.

In this experiment we only select 10 Hz frequency data for the reconstruction algorithm and underlying iterative minimization of the misfit functional (2.27). We simulate time-harmonic data using a Continuous Galerkin finite element method (CG). We use an approach and implementation similar to the one discussed in

⁶The code that was used, here, can be found at <https://team.inria.fr/magique3d/software/hou10ni/>

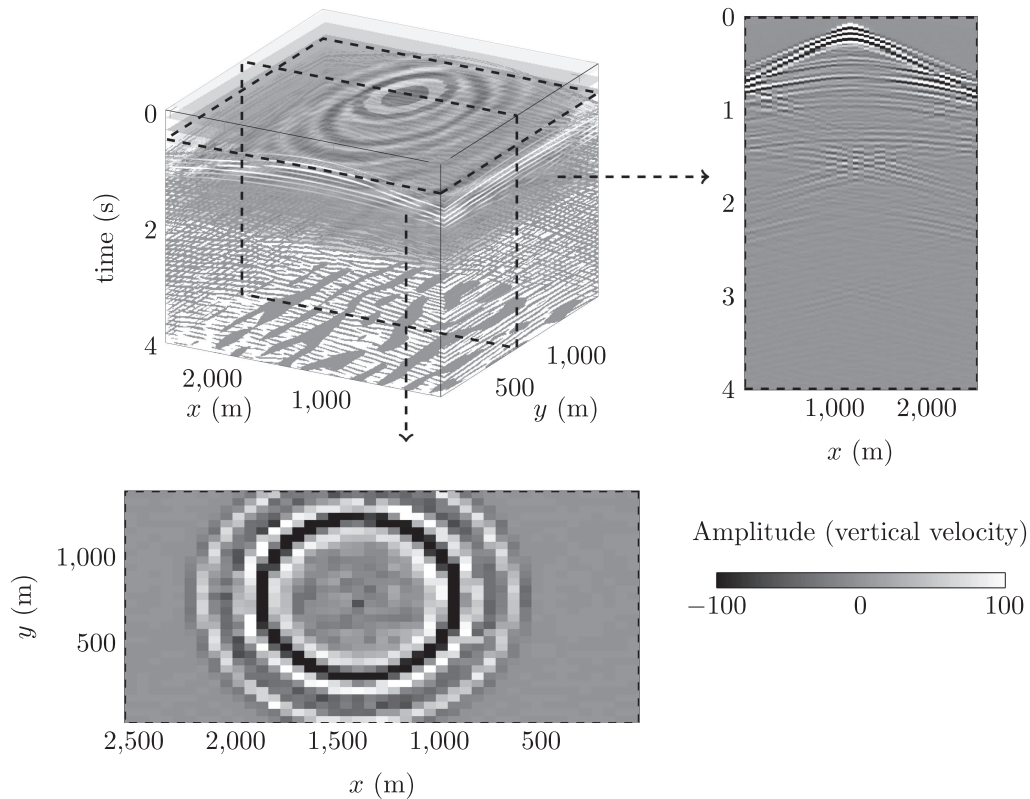


FIGURE 6. Three-dimensional time-domain vertical velocity trace for a centrally located source. Sections at fixed time $t = 0.5$ s and for a fixed line of receivers positioned in $y = 695$ m are respectively given at the bottom and right of the 3D visualization. The x and y axis correspond with the receivers map (*i.e.*, every 60 m for x and 45 m for y).

Shi *et al.* [35]. The relevant system of equations are solved with the direct structured solver MUMPS [5]. The numerical discretization introduces a tetrahedral representation of the model, which we illustrate in Figure 8. The choice of CG (instead of DG) is motivated by the memory cost of solving large linear system, which is a specificity of the harmonic case (DG is used for the time-domain discretization). In our experiments, we use order 3 polynomials to guarantee the accuracy (by taking at least four degrees of freedom per wavelength, according to the common heuristic). Note also that the mesh employed to generate the synthetic (time-domain) data differs from the one used for the inverse (harmonic) problem: the one to generate the data is refined to make sure we consider acutely the variations of the reference wave speed model.

For the reconstruction, the wave speed uses a piecewise linear representation, following (2.5b) and [3]. For the construction of the hierarchy of stable subspaces, the domain partition determining the piecewise linear representation of the wave speed is typically significantly coarser than the tetrahedral mesh. Hence, we define every subdomain, D_j , as the union of mesh elements, K_i , according to

$$D_j = \bigcup_{i=1}^{N_j} K_i, \quad (3.1)$$

where N_j denotes the number of mesh elements in D_j .

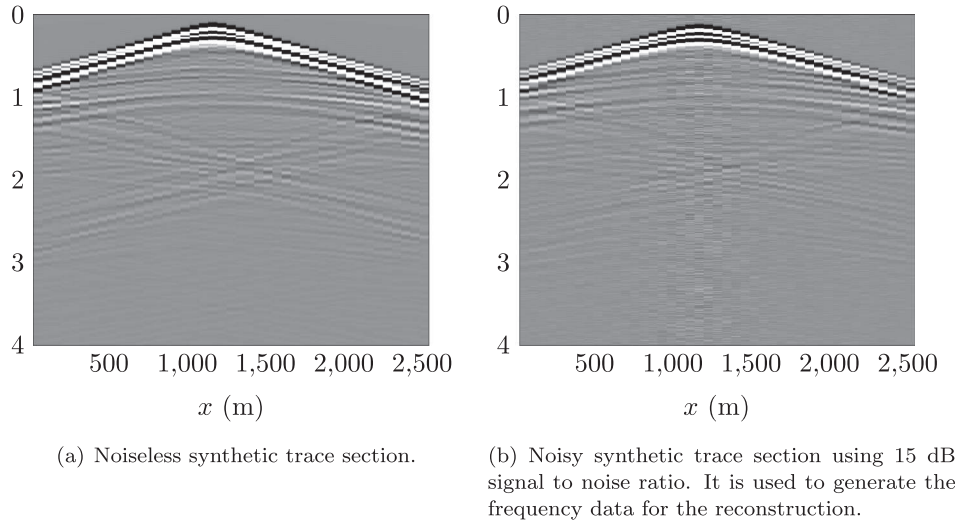


FIGURE 7. Comparison of noiseless and noisy data for a section of the pressure trace given Figure 5. The noise is independently generated for every receiver and every source in the acquisition.

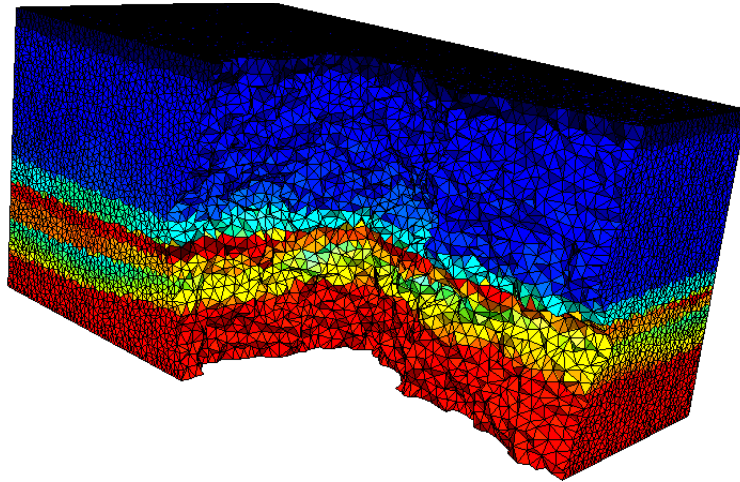


FIGURE 8. Illustration of the unstructured tetrahedral mesh of the model, using 433 979 tetrahedra.

To achieve the decomposition, we apply a structured decision where the maximal size of the subdomains is chosen in every direction to define the subspace. Then, piecewise linear functions are employed onto each generated subdomain to represent the model. We illustrate the effect of piecewise linear partitioning applied on the wave speed model in Figure 9, where the size of the subdomain is at most 400 m in the x and y directions, and 150 m in the z direction; this leads to a decomposition with $N = 224$ subdomains and 896 coefficients to represent the model. Inherent model error is introduced from those two levels of representation (the mesh and the partitioning). Because we do not know the subsurface geometry *a priori*, the piecewise linear partition relies

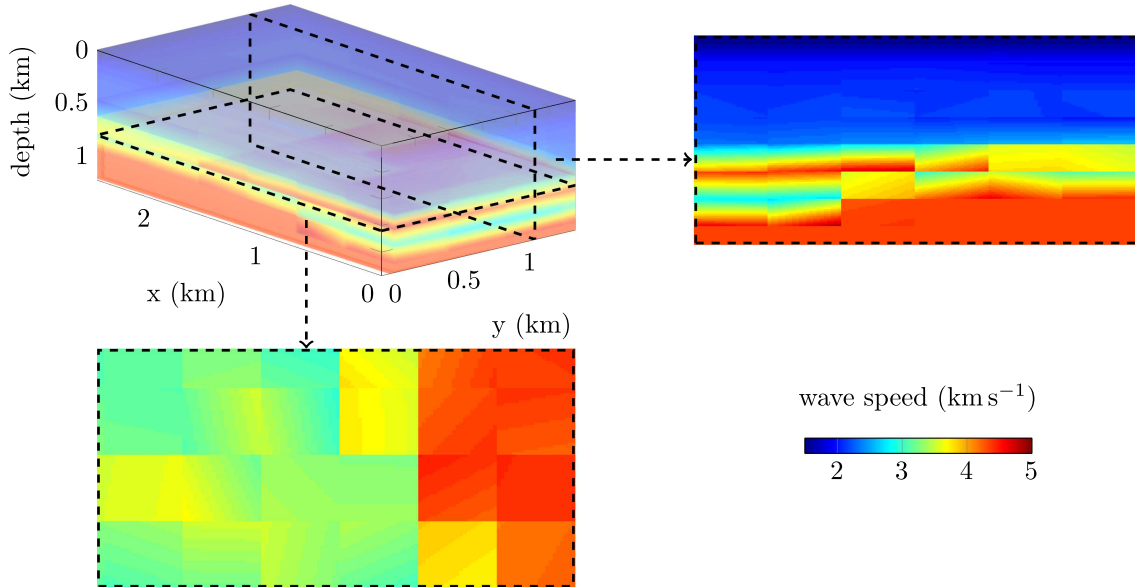


FIGURE 9. Illustration of three-dimensional partitioning of the wave speed model of Figure 3. The size of subdomains is limited to 400 m in the x and y directions and 150 m in the z direction. This leads to a decomposition of $N = 224$ subdomains where piecewise linear functions are used to represent the wave speed.

on the gradient of the misfit functional instead of the wave speed. Naturally, the more subdomains are taken, the more accurate can the representation be.

We proceed using single-frequency, 10 Hz data and a fixed domain partition. Exploiting the Lipschitz stability result obtained in Theorem 2.3 above, the Landweber iteration [18] provides a convergence analysis. The initial model needs to be within the radius of convergence. Algorithm 1 summarizes the steps of the reconstruction procedure.

Two parameters can decide of the termination of the procedure: if the number of iterations $n_{\text{iter}}^{(\max)}$ is reached, or if the cost function stagnates (criteria $\epsilon_{\mathcal{J}}$ and n_{ϵ}), see Algorithm 1. In the following experiments, we impose

$$\begin{aligned} n_{\text{iter}}^{(\min)} &= 50, & n_{\text{iter}}^{(\max)} &= 250, \\ n_{\epsilon} &= 10 & \text{and} & \quad \epsilon_{\mathcal{J}} = 0.01. \end{aligned} \tag{3.2}$$

Hence, the number of iterations is kept relatively high and the stagnation stops the procedure. More precisely with the given numbers, the algorithm stops if the difference in the misfit functional over the last ten iterations is less than 1%.

3.1. Single-frequency data

From the Cauchy data at 10 Hz, we carry out a reconstruction of the reference model starting from the smooth model depicted Figure 10. We encode the principal variation and appropriate order of magnitude of the wave speed in the initial model.

The model is partitioned in 1089 subdomains where piecewise linear functions are used to define the wave speed. This leads to a total number of unknowns of $4 \times 1089 = 4356$ (while the reference medium has 1 527 168 nodal values, see Fig. 3). The key, here, is the low-dimensional subspace used for regularization. The partition is adapted to the gradient computation *via* segmentation. We carry out 175 iterations for the two situations.

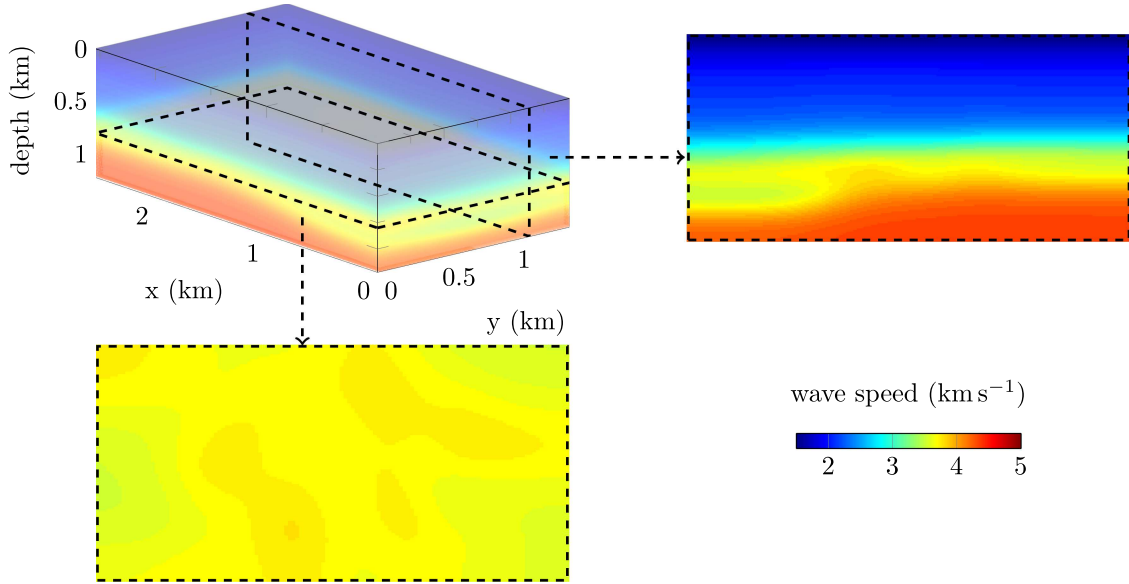


FIGURE 10. Three-dimensional representation, horizontal section at 800m depth and vertical section at $y = 1.125$ km of the initial wave speed model.

Algorithm 1: Algorithm for the reconstruction of subsurface parameters using piecewise linear model partition and Cauchy data.

Preliminary material:

- time-domain observation data at the receivers location,
- user prescribed frequency k and associated partition N ,
- user prescribed initial model c_1 and number of iterations $n_{\text{iter}}^{(\min)}$, $n_{\text{iter}}^{(\max)}$,
- user prescribed stagnation parameters $\epsilon_{\mathcal{J}}$ and n_{ϵ} .

Computation of the Fourier transform of the time-domain Cauchy traces at k .

Optimization loop for $j \in \{1, \dots, n_{\text{iter}}^{(\max)}\}$ do

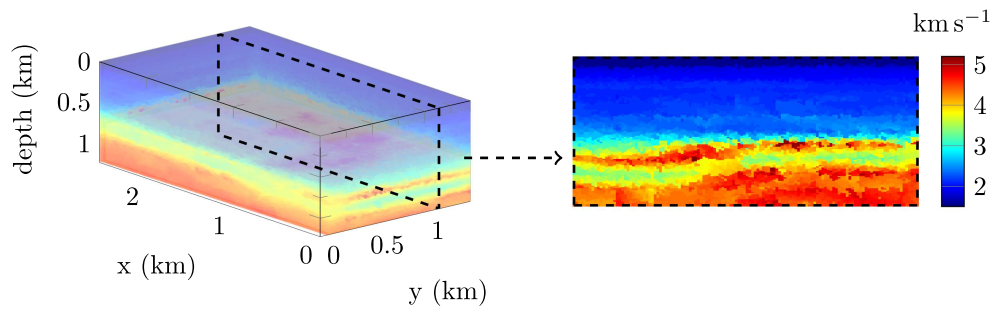
- solve the Helmholtz equation (2.16) at frequency k with wave speed c_j ;
- compute the misfit functional (2.27) from the simulation and observation data;
- compute the gradient of the misfit functional, $\nabla \mathcal{J}(c_j)$ with the adjoint-state method (see Sect. 2.4);
- compute the search direction, s_j , here we use the nonlinear conjugate gradient method with Polak–Ribière formula (cf. [28], Sect. 5.2);
- apply segmentation onto the search direction (see illustration on Fig. 9);
- compute the step length α with line search method (backtracking, cf. [28], Chap. 3);
- model update $c_{j+1} = c_j - \alpha s_j$;

if $j \geq n_{\text{iter}}^{(\min)}$ and $j > n_{\epsilon}$ then

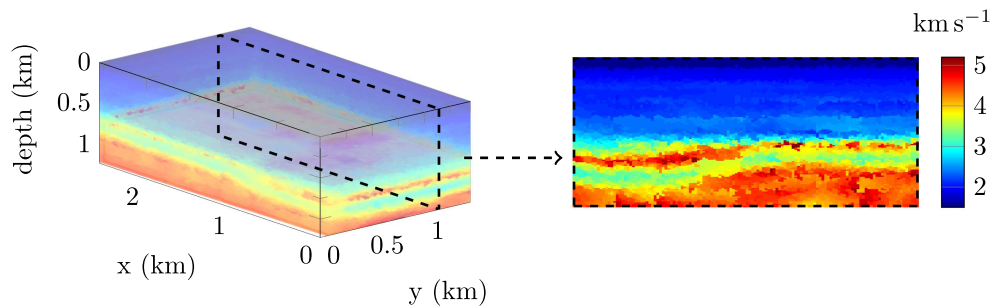
- compute the stagnation criterion $\epsilon = \frac{\mathcal{J}(c_{j-n_{\epsilon}}) - \mathcal{J}(c_j)}{\mathcal{J}(c_{j-n_{\epsilon}})}$;
- **if $(\epsilon < \epsilon_{\mathcal{J}})$: exit optimization loop.**

end

end



(a) Reconstruction where the sources in the acquisition are positioned in a three-dimensional area (see Fig. 4(b)); the relative L^2 difference with the reference model Figure 3 is $\mathcal{E}_{\text{rel}} = 0.083$.



(b) Reconstruction where the sources in the acquisition are limited on a plane (see Fig. 4(c)); the relative L^2 difference with the reference model Figure 3 is $\mathcal{E}_{\text{rel}} = 0.088$.

FIGURE 11. Three-dimensional representation and vertical section at $y = 1.125$ km of the reconstruction from 10 Hz Cauchy data after 175 iterations. The partition consists of $N = 1089$ subdomains, leading to a total number of unknowns of $4 \times 1089 = 4356$. The initial model is shown in Figure 10.

In Figure 11a, we show the reconstruction when the sources are positioned in a volume above the receivers (see Fig. 4b); in Figure 11b, we show the reconstruction when the sources are restricted on a two-dimensional plane (see Fig. 4c). To compare the accuracy of the reconstructions, we use the relative L^2 norm of the difference between the reference model and the reconstruction:

$$\text{relative } L^2 \text{ error} = \mathcal{E}_{\text{rel}} = \frac{\|c_{\dagger} - c_r\|}{\|c_{\dagger}\|}, \quad (3.3)$$

where c_{\dagger} is the reference model of Figure 3 and c_r the final reconstruction.

The two acquisitions provide accurate recoveries of the subsurface, with some improvement when the sources are positioned in a volume: the width of the increased wave speed layer and the deepest values of the wave speed are better retrieved with this type of acquisition. However, the reconstruction from sources limited on a plane is very close. In the subspace, we have drastically reduced the number of unknowns in the representation as compared with the original representation, namely to 0.3%. Nonetheless, the reconstruction captures the main features of the model, including the alternation of high and low values in the vertical direction and the resolution remains reasonable.

Remark 3.1 (Improved visualization with Gaussian filtering). The visualization of the reconstruction may suffer from the tetrahedral mesh employed for the numerical discretization. It is simple to improve the visualization by applying a smoothing filter onto the image. This can be done, for example, with the `imgaussfilt` function of

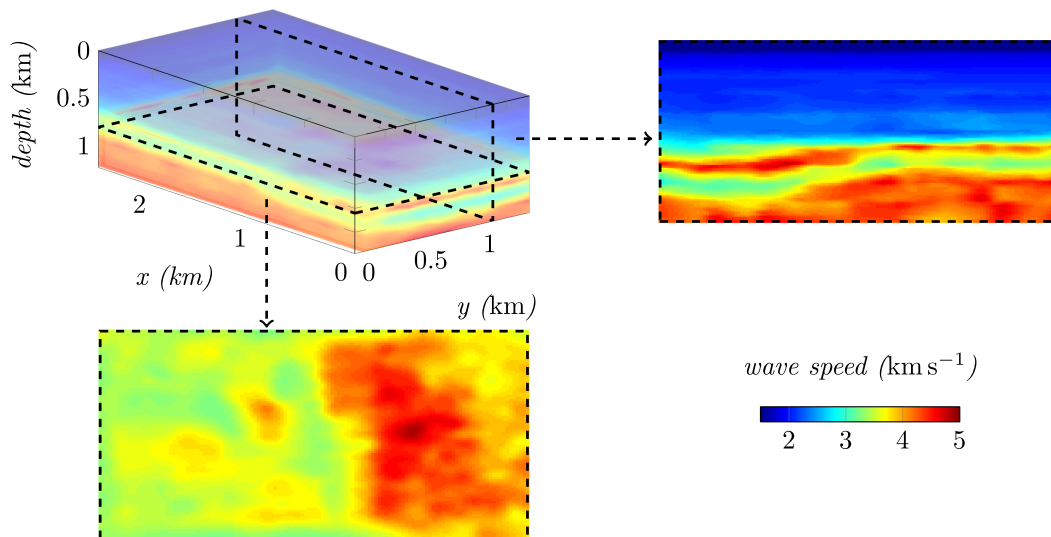


FIGURE 12. Gaussian filtering of the reconstruction obtained from 10 Hz Cauchy data (Fig. 11b), three-dimensional representation, horizontal section at 800 m depth and vertical section at $y = 1.125$ km. The relative L^2 difference with the reference model Figure 3 is $\mathcal{E}_{\text{rel}} = 0.079$.

MATLAB, which applies a Gaussian smoothing filter. In Figure 12, we show the resulting image when applied onto the reconstruction of Figure 11b. It allows a better identification of the recovered structures. Note that this procedure is done *a posteriori*, independently of the reconstruction algorithm, and is effortless.

In the following experiments, we only consider the case where sources are restricted on a plane, following the acquisition illustrated Figure 4c, for simplicity.

3.2. Single-frequency data, depth varying initial model

We repeat the experiment carried out in the previous subsection (with the two-dimensional sources lattice), but with a simplified initial model, see Figure 13. That is, here, the initial model only contains an indication of the average variation of wave speed in depth. The idea behind this experiment is to test the radius of convergence on the one hand, and the closeness of the true model and the best projection onto a low-dimensional stable subspace on the other hand.

In Figures 14 and 15, we present the result after 175 iterations using 10 Hz Cauchy data. As in the first experiment, we have 1089 subdomains in the partition and piecewise linear representations. Despite the lack of initial information we still retrieve the main features and appropriate contrasts in the wave speed. However, we lost accuracy as compared with the previous example (especially on the side), but the deep layer of low wave speed is well identified nonetheless.

In Figure 16, we show the evolution of the misfit functional with iterations on a logarithmic scale, where we compare with the previous experiment that used a smooth initial model (see Figs. 10 and 11b for the reconstruction). As expected, the first iteration using the smooth initial model provides a reduction in the misfit functional compared to the one-dimensional starting model. The decrease of the misfit functional is relatively fast for the initial iterations, especially when starting with the smooth model, and we observe a slow evolution after about 100 iterations in both configurations. Eventually, we observe the stagnation which stops the procedure. As indicated with the L^2 norm of the difference between the reference model (see Figs. 11b and 14), starting with the smooth model provides a better approximation.

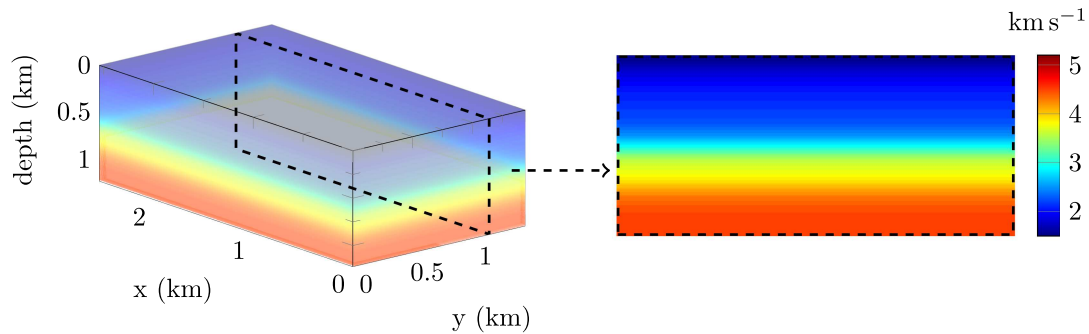


FIGURE 13. Three-dimensional representation and vertical section at $y = 1.125$ km of the initial model, which varies in depth only.

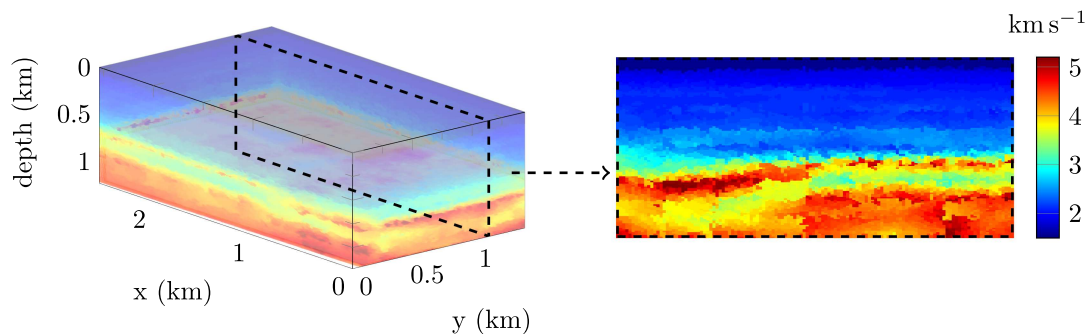


FIGURE 14. Three-dimensional representation and vertical section at $y = 1.125$ km of the reconstruction from 10 Hz Cauchy data after 175 iterations. The partition consists of $N = 1089$ subdomains, leading to a total number of unknowns of $4 \times 1089 = 4356$. The initial model varies in depth only and is shown in Figure 13. The relative L^2 difference with the reference model Figure 3 is $\mathcal{E}_{\text{rel}} = 0.119$.

Figures 17 and 18 compare the observed, initial and reconstruction data, using the full receivers map associated with a centrally located source. We see that the data from the recovered wave speed provide a pattern that is similar to the Fourier transform of the time-domain observations.

4. PERSPECTIVES ON INDEPENDENT LOCATIONS OF SOURCES IN THE DISCRETIZED SETTINGS

The misfit functional (2.27) defined for the Cauchy data has an interesting feature, because of the intuitive differentiation between acquisition sets for the observations and simulations. It is materialized by the double integral over $K_1 \times K_1$. The perspective is here to separate in, say, $K_2 \times K_1$.

In the usual context of minimization involving the direct difference between observations and simulations, such as the standard least squares, the setup for simulation is imposed by the field acquisition (source position and wavelet). Consequently, absence of knowledge leads to the failure of the algorithm. Here, providing this new misfit functional, we expect our iterative minimization algorithm to be free of those considerations, introducing extreme flexibility for the setup, where only the position of the receivers is required.

To illustrate the potential of the method, we design an experiment where the simulation sources differ from the observational ones. We consider a subsurface wave speed model where salt domes (objects having a large

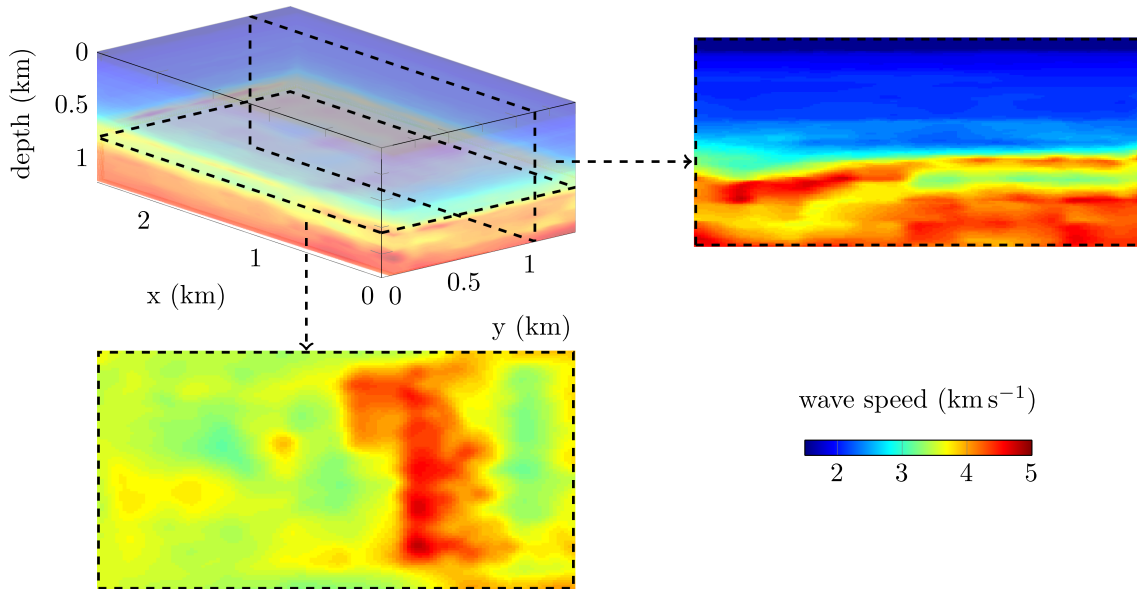


FIGURE 15. Gaussian filtering (see Rem. 3.1) of the reconstruction (Fig. 14) obtained from 10Hz Cauchy data starting with the initial model that varies in depth only (Fig. 13). Three-dimensional representation, horizontal section at 800 m depth and vertical section at $y = 1.125$ km. The relative L^2 difference with the reference model Figure 3 is $\mathcal{E}_{\text{rel}} = 0.108$.

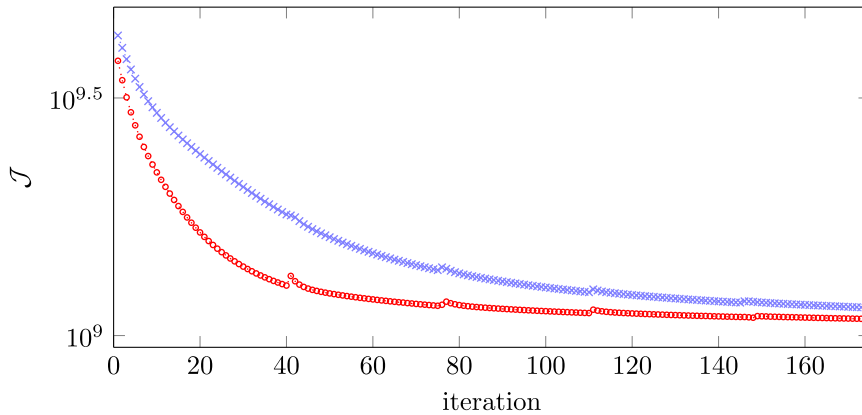


FIGURE 16. Evolution of the misfit functional with iterations, the blue crosses (-x-) correspond with the case where the initial model is one-dimensional (see Fig. 13 for the initial model and Fig. 14 for the associated reconstruction); the red circles (-o-) correspond with the case where the initial model is smooth (see Fig. 10 for the initial model and Fig. 11b for the associated reconstruction).

speed contrast) are present. The model is illustrated in Figure 19 and is of size $2.46 \times 1.56 \times 1.2$ km. It consists in a smooth background with contrasting objects having a wave speed of 4500 m s^{-1} . We assume that the density remains constant with $\rho = 1000 \text{ kg m}^{-3}$. This medium is very different in nature from the previous one, and salt domes are traditionally challenging in seismic exploration (especially where their presence is initially unknown).

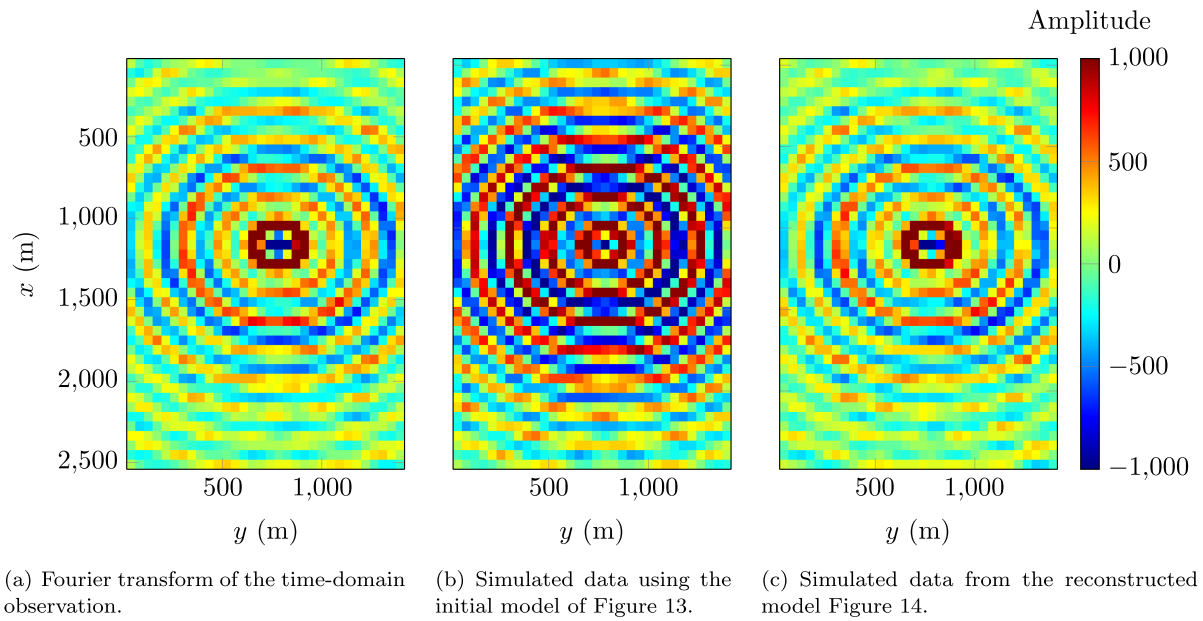


FIGURE 17. Comparison of the 10 Hz frequency pressure data captured at the receivers location for a centrally located source.

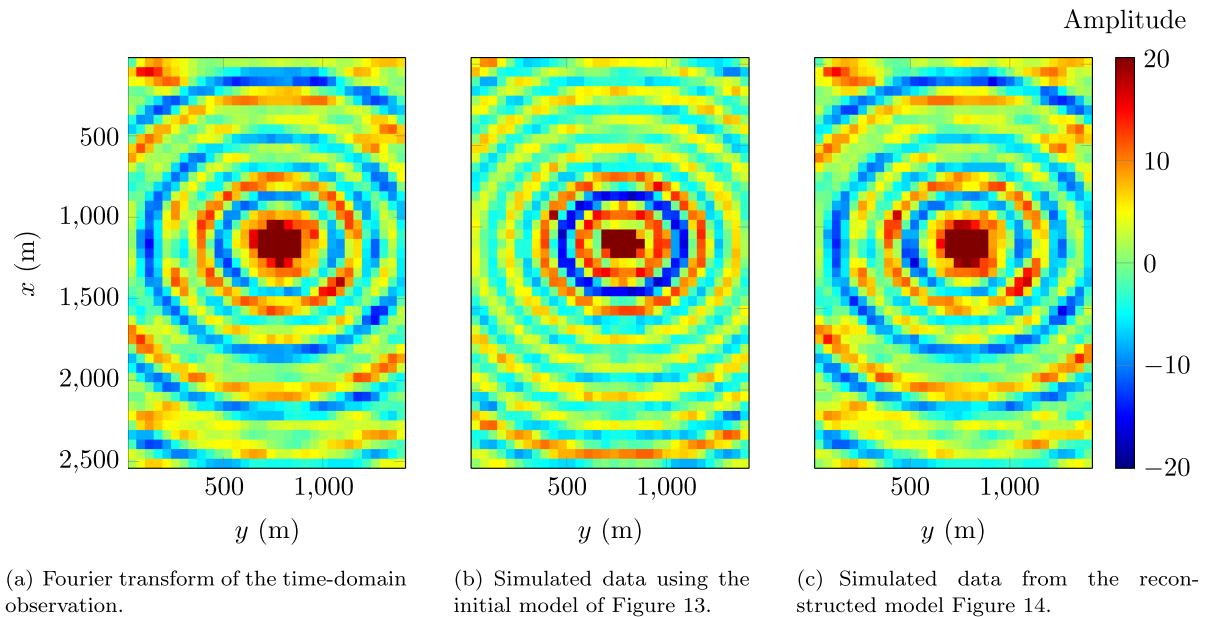


FIGURE 18. Comparison of the 10 Hz frequency vertical velocity data captured at the receivers location for a centrally located source.

Time-domain Cauchy data are obtained from this configuration with 1000 fixed receivers for each of the 96 sources. Here the sources are positioned on a plane, according to Figure 4c. The devices are located just underneath the surface, at a depth of 10 m for the source and 100 m for the receivers. We incorporate noise in

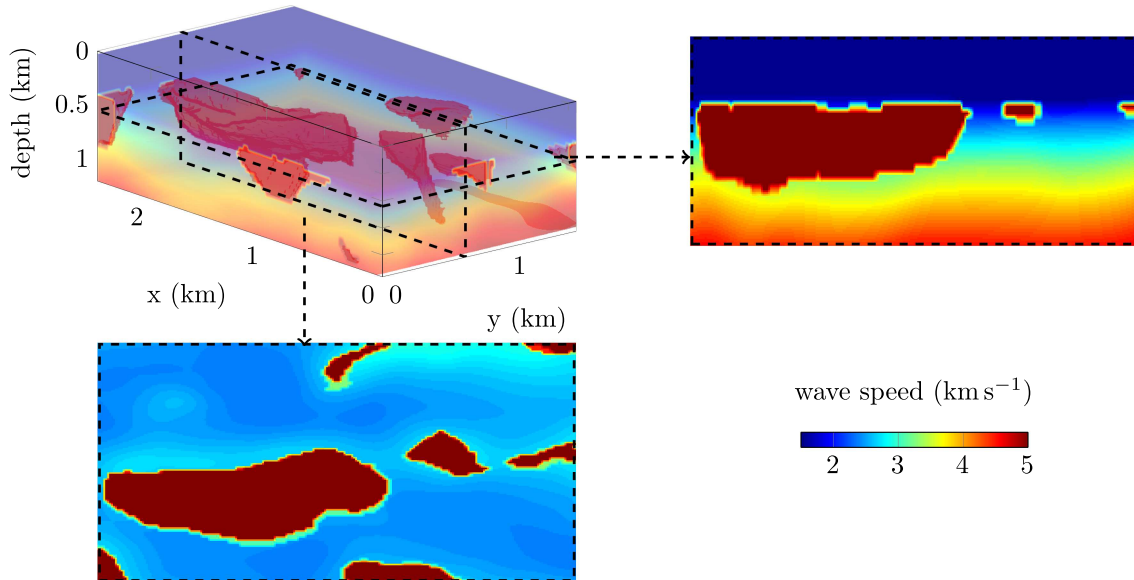


FIGURE 19. Three-dimensional representation, horizontal section at 550 m depth and vertical section at $y = 670$ m of the reference wave speed encompassing salt domes. The wave speed in the domes is a constant of 4500 m s^{-1} .

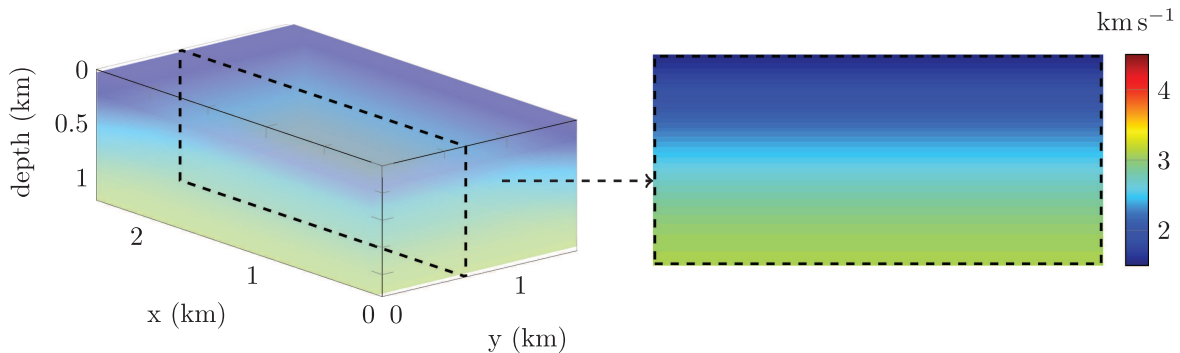


FIGURE 20. Three-dimensional representation and vertical section at $y = 670$ m of the initial model, which varies in depth only.

the time-domain data, with a signal-to-noise ratio of 15 dB, before we proceed to the Fourier transform to get the frequency-domain data. For the reconstruction we start with an initial model which only varies with depth, see Figure 20. We do not assume any contrasting objects in our initial guess, nor do we know the value for the background. For the reconstruction we only assume the knowledge of the uppermost water layer (up to 150 m depth).

To differentiate the simulation set of sources from the observation, we reduce their number, change their position and modify the source wavelet, see Table 1. We perform 100 iterations of the reconstruction Algorithm 1, with single frequency data at 4 Hz. The model representation is fixed with $N = 1280$ sub-domains where piecewise linear functions are employed, for a total of 5120 coefficients.

TABLE 1. Comparison of acquisition setups employed for the observations and simulations. The source wavelet also differs. For the reconstruction, single frequency data at 4 Hz and a fixed model partition of $N = 1280$ are used.

	Setup for measurements	Setup for simulations
Number of sources	96	60
Depth of the sources	10 m	20 m

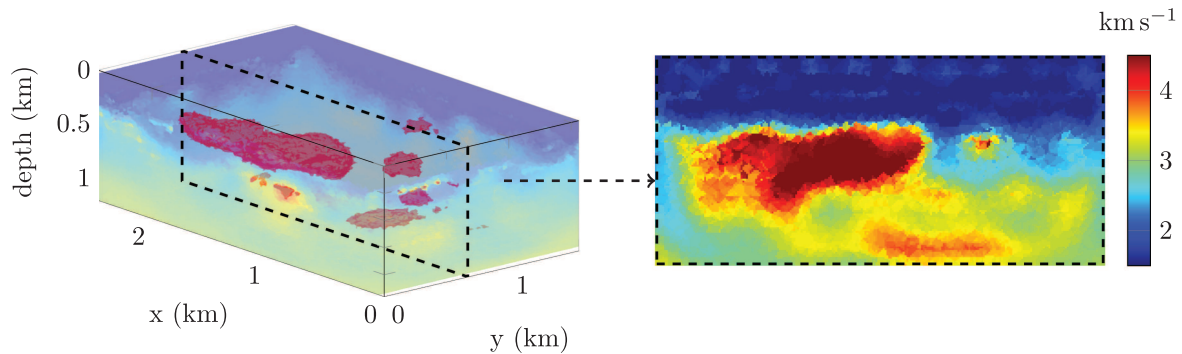


FIGURE 21. Three-dimensional representation and vertical section at $y = 670$ m of the reconstruction using 4 Hz Cauchy data after 100 iterations and where the simulation setup differs from the original measurement acquisition, see Table 1. The relative L^2 difference with the reference model Figure 19 is $\mathcal{E}_{\text{rel}} = 0.174$.

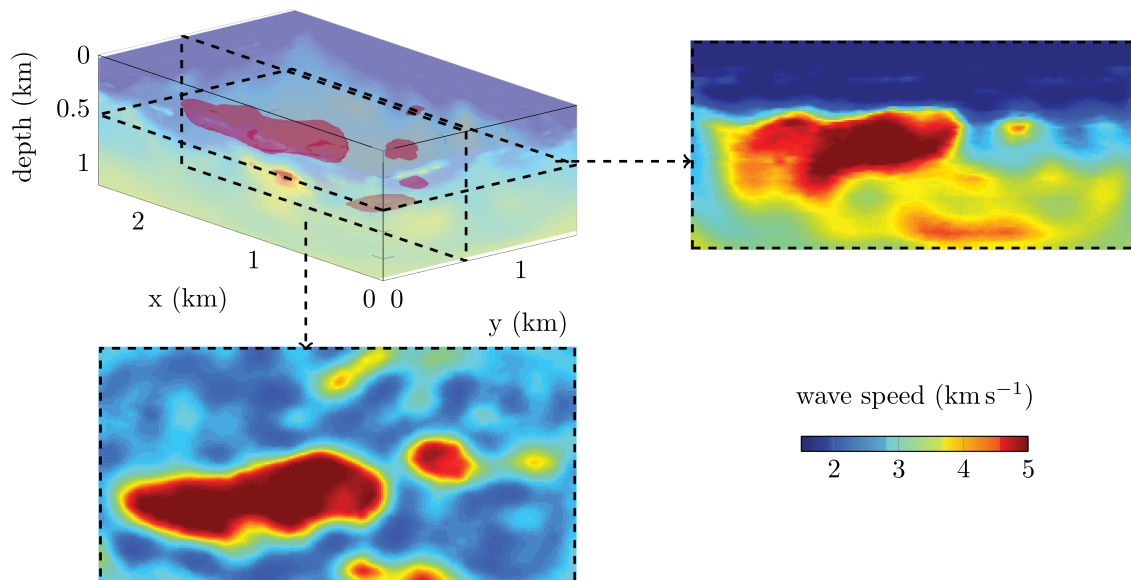


FIGURE 22. Gaussian filtering (see Rem. 3.1) of the reconstruction (Fig. 21) obtained using 100 iterations with 4 Hz Cauchy data and where the setup for simulation differs from the original measurement acquisition. Three-dimensional representation, horizontal section at 550 m depth and vertical section at $y = 670$ m. The relative L^2 difference with the reference model Figure 19 is $\mathcal{E}_{\text{rel}} = 0.172$.

The reconstruction after 100 iterations is shown Figures 21 and 22. Despite our initial guess having no information on the objects, the procedure is able to recover the main dome with the accurate value, and the shape of smaller domes (see the horizontal section in Fig. 22). The near boundary information seems missing, as well as the deepest model variation, due to limited illumination. However, as we started with a one dimensional guess and used single frequency data, the reconstruction is very satisfactory. Once again, the restricted number of unknowns does not prevent a good resolution.

We have changed the number of sources in the simulation compared to the observation, reducing the numerical cost accordingly, yet we make full use of the observed data from the benefit of our misfit functional defined for Cauchy data. Furthermore we do not need to know the position of the sources employed for observation, nor the source wavelet. The perspective of differentiating the observations and simulations acquisition sets is a promising application, and appears consistent with the results of this preliminary experiment. It would allow less prior on the observational environment, increasing the robustness of the procedure, without impacting the resolution of the reconstruction.

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