

THEORETICAL AND NUMERICAL STUDY OF THE CONVERGENCE OF LUENBERGER OBSERVERS FOR A LINEARIZED WATER WAVE MODEL

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Abstract. This paper investigates the convergence properties of Luenberger observers applied to a linearized water wave model. The study is motivated by the challenge of estimating wave dynamics when only partial free surface measurements are available. We identify fundamental obstructions to convergence, showing that the classical Luenberger observer fails to achieve full-state reconstruction due to challenges associated with mean value modes and high-frequency components. To overcome these limitations, we introduce modified observer schemes that incorporate frequency filtering and projection techniques. Our theoretical results are reinforced by numerical experiments that demonstrate the practical effectiveness of these observer-based estimation methods for water waves.

Key words. Linearized water wave equation, Luenberger observer, numerical simulations.

AMS subject classifications. 76B15, 93B53, 93C20

1. Introduction. The study of water waves has long been a cornerstone of fluid mechanics, with broad applications in oceanography, naval engineering, and coastal management. These wave dynamics are governed by the water wave equations (WWE), derived from the incompressible Euler equations, which describe the free surface's evolution and the underlying fluid's velocity potential. However, to simplify the complex nonlinearity of these equations, the linearized water wave equations (LWWE) are often employed. These equations, while capturing essential physical properties, allow for more tractable analytical and numerical analysis. They form the basis for the observer convergence analyses explored in this paper.

Recent research in the literature on fluid mechanics engineering science has focused on estimation and control techniques for wave systems, particularly in scenarios where direct measurements of the full state of the system are unavailable or impractical [9, 7, 10]. Luenberger observers, initially designed for linear systems, have been adapted to a range of complex applications, including fluid dynamics. The primary objective of this paper is to examine from a mathematical point of view the convergence properties of Luenberger observers when applied to the linearized water wave model, with partial data from the free surface (the case of full surface measurements has been studied in [36]). Although the existing literature offers a variety of methods for controlling and stabilizing water waves, there is a gap regarding the rigorous convergence analysis of such observers in the context of the LWWE with partial observation of the water surface.

Several works have investigated control and stabilization methods for water waves. Gagnon et al. [14] presented a Fredholm-type backstepping transformation to achieve rapid stabilization of the so-called capillary-gravity linearized water waves model, using spectral properties and controllability assumptions. We also refer to the contribution of Alazard et al. [1], who provided insights into the stabilization of water waves using pressure disturbances applied to the free surface. Further contributions in the field of stabilization include Su et al. [33], who investigated the stabilization of small-amplitude water waves using boundary controls. Note that the problem studied

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in this paper has similarities to that of damped wave equations. We refer to [31] for a comprehensive review of this problem.

While previous studies have concentrated on stabilization and controllability for both linear and nonlinear water wave models, our work provides a novel approach by focusing specifically on the convergence of Luenberger observers in a linearized setting. The contributions of this paper are twofold. First, we establish some obstructions for the convergence of Luenberger observers applied to the LWWE, together with sufficient conditions to recover some “partial” and quantitative convergence results, leveraging spectral methods and operator theory. Related works are [15, 4], where the authors also study Luenberger observers or back and forth nudging algorithms for infinite-dimensional linear systems that are not fully reconstructible, as in the present situation. However, the results given in these articles do not apply in our context. Second, we validate these theoretical results through numerical simulations, demonstrating the practical applicability of observer-based estimation for water wave systems.

This paper is organized as follows. Section 2 introduces the linearized water wave model. Section 3 sets the functional setting for our analysis. Section 4 delves into the main non-convergence results (Theorem 4.1 and Theorem 4.4) for the most natural Luenberger observer if we consider initial conditions in the natural energy space, or in a space of initial conditions with zero mean value. Section 5 aims to understand how to restore some partial and quantitative convergence theorems by changing the Luenberger observer (Theorem 5.1 and Theorem 5.2). Lastly, section 6 validates the theoretical analysis through numerical experiments and discusses the broader implications for observer design in fluid systems.

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2. The linearized water wave equation. This paper focuses on linearized water wave (LWWE) equations. These equations are derived from a master model, namely the water wave equations (WWE), itself derived from the incompressible Euler equations. These derivations can be found in [11] or [23]. Our geometric setting is summarized in Figure 2.0. We place ourselves in a domain Ω that describes the body of our fluid:

$$\Omega = \{(x, y) \in \mathbb{R}^2 \mid -d \leq y \leq \eta(t, x)\}.$$

The two main functions of interest are: $\psi(t, x, y)$ the velocity potential at points (x, y) and time t , and $\eta(t, x)$, the free surface elevation at point x and time t . $-d$ is a flat bottom. As described in [11] or [23] for water wave with a horizontal bottom and in [6] if we add the periodicity, ψ is harmonic in space:

$$(2.1) \quad \Delta\psi(t, x, y) = 0, \quad (x, y) \in \Omega, \quad t \in [0, +\infty).$$

Moreover, ψ is subject to the following boundary conditions:

$$(2.2) \quad \begin{aligned} \partial_y \psi(t, x, y) &= 0, & \text{on } y = -d, \\ \partial_t \psi(t, x, y) &= -g\eta(t, x), & \text{on } y = 0, \\ \partial_t \eta(t, x) &= \partial_y \psi(t, x, y), & \text{on } y = 0, \\ \psi(t, x, y) &= \psi(t, x + L, y), & (t, x, y) \text{ in } [0, +\infty) \times \Omega. \end{aligned}$$

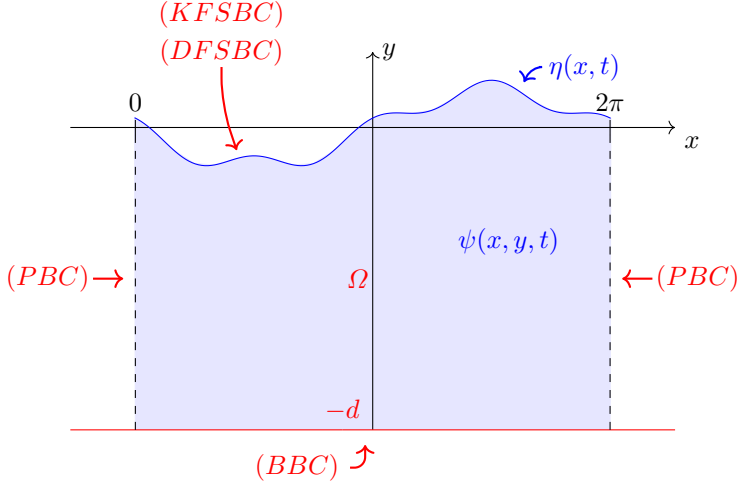


FIGURE 2.1. Boundary problem for periodic water wave with flat bottom.

85 The first boundary condition is referred to as the *bottom boundary condition*
86 (*BBC*), the second as the *dynamic free surface boundary condition* (*DFSBC*), the
87 third as the *kinematic free surface boundary condition* (*KFSBC*), and the last as
88 the *periodic boundary condition* (*PBC*). Readers are again referred to [6], [23] or
89 [11, Section 2.3] to see how the linearization around the rest state $y = 0$ of the
90 Bernoulli equation leads to those conditions. Moreover, for computational simplicity,
91 we consider here x being on the torus $\mathbb{T} = \mathbb{R}/2\pi\mathbb{Z}$, *i.e.* choosing $L = 2\pi$. Under
92 all these assumptions and after applying the boundary conditions, it is shown [11,
93 2.3] that the velocity potential $\psi(t, x, y)$ can be expressed simply at the surface
94 as a function $\phi(t, x) = \psi(x, 0, t)$, and reconstructed throughout the whole domain Ω .
95 This results in the following system of equations which is now on variables ϕ and η ,
96 respectively, the velocity potential at the surface and the free surface displacement:

$$97 \quad (2.3) \quad \begin{cases} \partial_t \phi(t, x) &= -g\eta(t, x), \\ \partial_t \eta(t, x) &= \mathcal{G}(\phi)(t, x), \quad (t, x) \in [0, +\infty) \times \mathbb{T}, \\ \phi(0, x) = \phi^0(x), \quad \eta(0, x) &= \eta^0(x), \end{cases}$$

98 where the operator $\mathcal{G}$ can be viewed as a Dirichlet-to-Neumann map that does the
99 link between the velocity potential throughout the whole domain ψ , and the velocity
100 potential at the surface ϕ . This operator is such that $\mathcal{G} : \phi \mapsto \partial_y \psi(t, x, y)|_{y=0}$. Thus,
101 following [11, Section 2.3], the operator $\mathcal{G}$ is in Fourier space

$$102 \quad (2.4) \quad \mathcal{F}[\mathcal{G}(\phi)](n) = |n| \tanh(d|n|) \mathcal{F}[\phi](n), \quad n \in \mathbb{Z}.$$

103 where $\mathcal{F}[\cdot]$ is the Fourier transform on $\mathbb{T}$. From (2.4), it is easy to obtain that functions
104 of the form:

$$105 \quad (2.5) \quad (\phi, \eta)^t = (\phi_0 e^{i(nx - \omega_n t)}, \eta_0 e^{i(nx - \omega_n t)})^t$$

106 with ω_n the angular frequency, are solutions of (2.3), provided that

$$107 \quad (2.6) \quad i\omega_n \phi_0 = g\eta_0,$$

108 and the following (dispersion) relation holds:

$$109 \quad (2.7) \quad (\omega_n)^2 = g|n| \tanh(d|n|).$$

110 **3. Functional setting.** We now introduce an appropriate functional setting.
 111 Let us introduce the spaces

$$112 \quad L_p^2 = L^2(\mathbb{T}) = \left\{ f = \sum_{n \in \mathbb{Z}} f_n e^{inx} \text{ such that } \|f\|_{L_p^2}^2 := \sum_{n \in \mathbb{Z}} |f_n|^2 < +\infty \right\},$$

$$113 \quad H_p^{1/2} = \left\{ f = \sum_{n \in \mathbb{Z}} f_n e^{inx} \in L_p^2 \text{ such that } \|f\|_{H_p^{1/2}}^2 := |f_0|^2 + \sum_{n \in \mathbb{Z}^*} |n| |f_n|^2 < +\infty \right\},$$

$$114 \quad H_p^1 = \left\{ f = \sum_{n \in \mathbb{Z}} f_n e^{inx} \in L_p^2 \text{ such that } \|f\|_{H_p^1}^2 := |f_0|^2 + \sum_{n \in \mathbb{Z}^*} n^2 |f_n|^2 < +\infty \right\}.$$

117 Each of these spaces is endowed with the scalar product naturally associated with
 118 the norm appearing in the definition. With these scalar products, they are Hilbert
 119 spaces. Note that, notably, by the Plancherel theorem,

$$120 \quad (3.1) \quad \|f\|_{L_p^2}^2 = \frac{1}{2\pi} \int_{\mathbb{T}} |f|^2, \forall f \in L_p^2.$$

121 Thanks to (2.4), the operator $\mathcal{G}$ can be expressed as follows:

$$122 \quad (3.2) \quad \forall f = \sum_{n \in \mathbb{Z}} f_n e^{inx} \in H_p^1, \mathcal{G}f(x) = \sum_{n \in \mathbb{Z}^*} |n| \tanh(d|n|) f_n e^{inx} (\in L_p^2).$$

123 Now, consider the operator A defined as:

$$124 \quad (3.3) \quad A := \begin{pmatrix} 0 & -gI \\ \mathcal{G} & 0 \end{pmatrix},$$

125 with state space $H = H_p^{1/2} \times L_p^2$ and domain $D(A) = H_p^1 \times H_p^{1/2}$, endowed respectively
 126 with the Hilbertian norms (for which they are Hilbert spaces)

$$127 \quad \|(f, g)\|_H^2 = \|f\|_{H_p^{1/2}}^2 + \|g\|_{L_p^2}^2, \|(f, g)\|_{D(A)}^2 = \|f\|_{H_p^1}^2 + \|g\|_{H_p^{1/2}}^2.$$

128 In this setting, the linearized water wave equations for periodic waves with an hori-
 129 zontal bottom given in (2.3) can be written as:

$$130 \quad (3.4) \quad \begin{cases} \partial_t y(t, x) &= Ay(t, x), \\ y(0, x) &= y^0(x), \end{cases} (t, x) \in [0, +\infty) \times \mathbb{T},$$

131 where $y = (\phi, \eta)^t$ and $y^0 = (\phi^0, \eta^0)^t$.

132 **3.1. Spectral analysis of $\mathcal{G}$ and A .** With the expression (3.2), it is easy to
 133 see that the distinct eigenvalues of the self-adjoint operator $\mathcal{G}$ are given by the ei-
 134 genvalue 0 (associated with the constant eigenfunction 1), and by the eigenvalues
 135 $\mu_k = |k| \tanh(d|k|)$, $k \in \mathbb{N}^*$, which are of multiplicity 2, with the associated orthogo-
 136 nal basis of eigenfunctions given by e^{ikx} and e^{-ikx} .

Contrary to many usual situations, let us emphasize that 0 is an eigenvalue of A , so that A is not a positive definite operator, but only a semi-definite operator. Hence, the usual theory to pass from first-order to second-order operators (see *e.g.* [34, Section 6.8]) cannot be applied, and notably, A is *not* a skew-adjoint operator. This will be an additional difficulty, taking into account that many results in the literature of Luenberger observers (or stabilization) are restricted to this case (see notably various results in [24, 15, 4, 5]).

However, it is quite easy to discover that A is “almost” skew-adjoint and to derive an appropriate spectral decomposition. The distinct eigenvalues of A are given by the two families

$$\lambda_k^+ = i\sqrt{g\mu_k} = i\sqrt{gk \tanh(dk)}, \quad \lambda_k^- = -i\sqrt{g\mu_k} = -i\sqrt{gk \tanh(dk)}, \quad k \in \mathbb{N}^*,$$

each different eigenvalue being of algebraic and geometric multiplicity 2, together with the eigenvalue 0, which is of algebraic multiplicity 2 and geometric multiplicity 1 (which means that we have a generalized eigenfunction associated with the eigenvalue 0). For $|n| = k$ with $k \neq 0$, a basis of orthonormal eigenfunctions associated with the corresponding eigenspace is given by

$$(3.5) \quad \nu_n^+ = \frac{1}{\sqrt{2|n|g}} \begin{pmatrix} i\sqrt{g}e^{inx} \\ \sqrt{|n|}e^{inx} \end{pmatrix}, \quad \nu_n^- = \frac{1}{\sqrt{2|n|g}} \begin{pmatrix} -i\sqrt{g}e^{inx} \\ \sqrt{|n|}e^{inx} \end{pmatrix}.$$

A normalized eigenfunction associated with 0 is given by

$$(3.6) \quad \nu_0^+ = \frac{1}{\sqrt{2}} \begin{pmatrix} 1 \\ 0 \end{pmatrix}.$$

and an (orthogonal) generalized and normalized eigenfunction is given by

$$(3.7) \quad \nu_0^- = \frac{1}{\sqrt{2}} \begin{pmatrix} 0 \\ 1 \end{pmatrix}.$$

The family of usual and generalized eigenfunctions forms a Hilbert basis in L_p^2 .

3.2. Well-posedness and conserved quantities. Let us introduce some additional notation, that will be useful for proving the well-posedness of (2.3). We define

$$L_{p,0}^2 = \left\{ f \in L_p^2 \text{ such that } \int_{\mathbb{T}} f = 0 \right\} = \left\{ \sum_{n \in \mathbb{Z}} f_n e^{inx} \in L_p^2 \text{ such that } f_0 = 0 \right\},$$

$$H_{p,0}^{1/2} = \left\{ f = \sum_{n \in \mathbb{Z}^*} f_n e^{inx} \in L_{p,0}^2 \text{ such that } \|f\|_{H_{p,0}^{1/2}}^2 := \sum_{n \in \mathbb{Z}^*} |n| |f_n|^2 < +\infty \right\},$$

$$H_{p,0}^1 = \left\{ f = \sum_{n \in \mathbb{Z}} f_n e^{inx} \in L_{p,0}^2 \text{ such that } \|f\|_{H_{p,0}^1}^2 := \sum_{n \in \mathbb{Z}^*} n^2 |f_n|^2 < +\infty \right\}.$$

These spaces are endowed with the natural scalar product and norms induced by their definitions, which we denote respectively by $\|\cdot\|_{L_{p,0}^2}, \|\cdot\|_{H_{p,0}^{1/2}}, \|\cdot\|_{H_{p,0}^1}$ (notice that we can extend these norms as semi-norms respectively on $L_p^2, H_p^{1/2}, H_p^1$, which corresponds to forgetting the mode 0, *i.e.* the mean, of the functions under study,

171 and that these norms are nothing else than the restrictions of the norms $\|\cdot\|_{L_p^2}, \|\cdot\|_{H_p^{1/2}}, \|\cdot\|_{H_p^1}$ on $L_{p,0}^2, H_{p,0}^{1/2}, H_{p,0}^1$.

172 $\|\cdot\|_{H_p^{1/2}}, \|\cdot\|_{H_p^1}$ on $L_{p,0}^2, H_{p,0}^{1/2}, H_{p,0}^1$.

173 Following the proof of [36, Theorem 1], we obtain that A , with domain

$$174 \quad \mathcal{D}_0(A) = H_{p,0}^1 \times H_{p,0}^{1/2},$$

175 is the generator of a C^0 -semigroup of operators on $H_0 := H_{p,0}^{1/2} \times L_{p,0}^2$ endowed with
 176 the natural scalar product, so that for any $(\phi^0, \eta^0) \in H_0$, there exists a unique solution
 177 $(\phi, \eta) \in C^0([0, +\infty), H_0)$ (this is exactly in the setting of [36], in this case). Moreover,
 178 it is easy to check that in the space H_0 , $D_0(A^*) = D_0(A)$ and $A^* = -A$, so that A is
 179 in fact a generator of a C^0 -group of operators on H_0 .

180 Our next goal is to understand how to obtain a well-posedness result on the whole
 181 state space H . This is the purpose of the following proposition.

182 **PROPOSITION 3.1.** *A is the generator of a C^0 -group of operators on H , so that*
 183 *for any $(\phi^0, \eta^0) \in H$, there exists a unique solution $(\phi, \eta) \in C^0([0, +\infty), H)$ to (2.3).*

184 *Proof.* We write $H_p^{1/2} \times L_p^2 = H_0 \oplus_\perp E_0$, where H_0 is in fact the closure of
 185 the nonzero eigenspaces of A and E_0 is the generalized eigenspace associated with 0.
 186 Associated with this decomposition, we have a natural decomposition of $A = A_0 + A_1$,
 187 where

$$188 \quad A_{0|H_0} = A, A_{0|E_0} = 0 \text{ and } A_{1|H_0} = 0, A_{1|E_0} = A.$$

189 Hence, by the previous discussion, it follows that A_0 is the generator of a strongly
 190 continuous semigroup on H . Moreover, A_1 is clearly a bounded operator, since E_0
 191 is a finite-dimensional space. We deduce (see *e.g.* [34, Theorem 2.11.2]) that A
 192 is indeed the generator of a C^0 -semigroup of operators on H . Moreover, the same
 193 analysis applies if we replace A by $-A$ (because A_0 is skew-adjoint, so $-A_0$ is also
 194 the generator of a strongly continuous semigroup on H , and $-A_1$ is still bounded),
 195 so $-A$ is also the generator of a strongly continuous semigroup on H . We deduce by
 196 applying [34, Proposition 2.7.8] that A is the generator of a strongly continuous group
 197 on H . $\square$

198 For a solution $y = (\phi, \eta)$ of (3.4), we introduce the following energy

$$199 \quad (3.8) \quad E(y, t) = \frac{1}{2} \left(g \int_{\mathbb{T}} (\eta(t, x))^2 dx + \int_{\mathbb{T}} (\sqrt{\mathcal{G}} \phi(t, x))^2 dx \right).$$

200 Formally differentiating this expression, using (3.4), and the fact that $\mathcal{G}$ is self-adjoint,
 201 we deduce that

$$\begin{aligned} \frac{d}{dt} E(y, t) &= g \int_{\mathbb{T}} \eta(t, x) \partial_t \eta(t, x) dx + \int_{\mathbb{T}} \sqrt{\mathcal{G}} \phi(t, x) \partial_t (\sqrt{\mathcal{G}} \phi)(t, x) dx \\ &= g \int_{\mathbb{T}} \eta(t, x) \mathcal{G} \phi(t, x) dx + \int_{\mathbb{T}} \sqrt{\mathcal{G}} \phi(t, x) \sqrt{\mathcal{G}} (\partial_t \phi)(t, x) dx \\ &= g \int_{\mathbb{T}} \eta(t, x) \mathcal{G} \phi(t, x) dx + \int_{\mathbb{T}} \mathcal{G} \phi(t, x) (\partial_t \phi)(t, x) dx \\ &= g \int_{\mathbb{T}} \eta(t, x) \mathcal{G} \phi(t, x) dx + \int_{\mathbb{T}} \mathcal{G} \phi(t, x) (-g \eta(t, x)) dx \\ &= 0, \end{aligned}$$

203 so that $E(y, t)$ is conserved:

$$204 \quad (3.9) \quad E(y, t) = \frac{1}{2} \left(g \int_{\mathbb{T}} (\eta^0)^2 + \int_{\mathbb{T}} (\mathcal{G}\phi^0)^2 \right).$$

205 These computations can be made rigorous by taking initial conditions in $D(A)$ and
206 using an easy density argument.

207 Let us give some invariant quantities that are inherent to (3.4). By the definition
208 of $\mathcal{G}$ given in (3.2), for any $f \in H^{1/2}$, we have $\int_{\mathbb{T}} \mathcal{G}f = 0$. Notably, using the second
209 equation of (3.4) and integrating in space on $\mathbb{T}$, formally, any solution (ϕ, η) is such
210 that

$$211 \quad \partial_t \left(\int_{\mathbb{T}} \eta(t, \cdot) \right) = \int_{\mathbb{T}} \mathcal{G}\phi(t, \cdot) = 0,$$

212 so that $\int_{\mathbb{T}} \eta(t, \cdot)$ remains constant over time:

$$213 \quad (3.10) \quad \forall t \geq 0, \int_{\mathbb{T}} \eta(t, \cdot) = \int_{\mathbb{T}} \eta^0.$$

214 Hence, the first equation of (3.4) integrated in space on $\mathbb{T}$ also gives that

$$215 \quad \partial_t \left(\int_{\mathbb{T}} \phi(t, \cdot) \right) = -g \int_{\mathbb{T}} \eta^0,$$

216 so that

$$217 \quad (3.11) \quad \forall t \geq 0, \int_{\mathbb{T}} \phi(t, \cdot) = \int_{\mathbb{T}} \phi^0 - gt \int_{\mathbb{T}} \eta^0.$$

218 These quantities are fully determined as soon as the initial conditions are known.
219 Conversely, knowing the mean value of ϕ and η at any time $t \geq 0$ is enough to recover
220 the zero mode of η^0 and ϕ^0 . Hence, if we assume that for some extra reason, we are
221 able to reconstruct or guess what the mean value of $\eta(t, \cdot)$ and $\phi(t, \cdot)$ are, then, it is
222 possible to reconstruct the zero mode of η^0 and ϕ^0 . Here also, these computations
223 can be made rigorous by taking initial conditions in $D(A)$ and using an easy density
224 argument.

225 The core idea of this paper is that, unlike in [36], we place ourselves in a setting
226 where we only get a partial measurement of the surface $\eta(t, x)$: $y(t, x) = \mathbf{1}_{\omega}(x)\eta(t, x)$,
227 where ω is an open subset of $\mathbb{T}$ that is distinct from $\mathbb{T}$. However, contrarily to the
228 result given in [36], this partial measurement is not enough to reconstruct the solution
229 of (3.4), as it is proved later on, and we need to find another strategy to recover at
230 least low frequencies different from 0. The first Luenberger-type observer we will
231 study is the following natural one:

$$232 \quad (3.12) \quad \begin{cases} \partial_t \hat{\phi}(t, x) &= -g\hat{\eta}(t, x), \\ \partial_t \hat{\eta}(t, x) &= \mathcal{G}\hat{\phi}(t, x) - \gamma \mathbf{1}_{\omega}(x)(\eta(t, x) - \hat{\eta}(t, x)), \\ \hat{\phi}(0, x) &= \hat{\phi}^0(x), \quad \hat{\eta}(0, x) = \hat{\eta}^0(x), \end{cases}$$

233 where $\gamma > 0$ is the correction gain, and $(\hat{\phi}, \hat{\eta})^T$ is the observer trajectory of the state
234 $(\phi, \eta)^T$.

235 Then, $(\phi_{er}, \eta_{er}) := (\phi - \hat{\phi}, \eta - \hat{\eta})$ is the solution to

$$236 \quad (3.13) \quad \begin{cases} \partial_t \phi_{er}(t, x) &= -g\eta_{er}(t, x), \\ \partial_t \eta_{er}(t, x) &= \mathcal{G}\phi_{er}(t, x) - \gamma (\mathbf{1}_{\omega}(x)\eta_{er}(t, x)), \\ \phi_{er}(0, x) &= \phi_{er}^0(x), \quad \eta_{er}(0, x) = \eta_{er}^0(x), \end{cases}$$

237 where $\phi_{er}^0(x) = \phi^0(x) - \hat{\phi}^0(x)$ and $\eta_{er}^0(x) = \eta^0(x) - \hat{\eta}^0(x)$.

238 To conclude, let us introduce the output and observer gain operators. We consider
 239 the output space Y to be the same as the state space H . They are given respectively
 240 by

$$241 \quad (3.14) \quad C = \begin{pmatrix} 0 & 0 \\ 0 & 1_\omega \end{pmatrix} \in \mathcal{L}_c(H) \text{ and } L = \gamma I_d \in \mathcal{L}_c(H),$$

242 so that (3.13) can be also abstractly written as

$$243 \quad (3.15) \quad \begin{cases} \partial_t \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix} = (A - LC) \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix}, \\ \begin{pmatrix} \phi_{er}(0, x) \\ \eta_{er}(0, x) \end{pmatrix} = \begin{pmatrix} \phi_{er}^0(x) \\ \eta_{er}^0(x) \end{pmatrix}. \end{cases}$$

244 4. Non-reconstruction results for the localized Luenberger observer on 245 the whole state space.

246 **4.1. Obstruction coming from the mean value.** Let us prove the following
 247 easy fact.

248 PROPOSITION 4.1. (3.13) is not asymptotically stable: there exists $y^0 \in H$ such
 249 that the corresponding solution y to (3.13) does not satisfy $\|y(t)\|_H \rightarrow 0$ as $t \rightarrow +\infty$.

250 *Proof.* This is just a question of noting that ν_0^+ defined in (3.6) is still an eigen-
 251 function of the operator $A - LC$, associated with the eigenvalue 0. Indeed, we have

$$252 \quad (A - LC)\nu_0^+ = A\nu_0^+ - \gamma L\nu_0^+ = \begin{pmatrix} 0 \\ 0 \end{pmatrix}.$$

253 Hence, the corresponding solution to (3.13) is constant in time and given by

$$254 \quad (\phi_{er}(t, x), \eta_{er}(t, x)) = \nu_0^+.$$

255 Hence, for such a solution, $\|y(t)\|_H$ is a nonzero constant and cannot converge to 0
 256 as $t \rightarrow +\infty$. $\square$

257 REMARK 4.2. Proposition 4.1 also holds if we replace $L = \gamma I_d$ by any $L \in \mathcal{L}_c(H)$,
 258 or even for any $L : H \rightarrow H$ for which the equation (3.15) makes sense, as is readily
 259 seen from the proof.

260 According to Theorem 4.1, it is tempting to try to stabilize all frequencies except the
 261 0 one, corresponding to the mean value.

262 **4.2. Obstruction coming from the high frequencies for the original Lu-
 263 enberger observer.** A key point in proving that the Luenberger observer does not
 264 converge exponentially even if we do not wish to reconstruct the mean value is the
 265 following result, which might be of independent interest.

266 LEMMA 4.3. Let $\mathcal{H}$ and $\mathcal{U}$ be two Hilbert spaces. Assume that $\mathcal{A}$ is the generator
 267 of a C^0 -group of operators on $\mathcal{H}$, and that $\mathcal{B} \in \mathcal{L}_c(\mathcal{U}, \mathcal{H})$. Assume that $\mathcal{P} \in \mathcal{L}_c(\mathcal{H})$ is
 268 a projection onto the closed subspace $\mathcal{F} = \mathcal{P}(\mathcal{H})$, and that $\mathcal{P}$ commutes with $\mathcal{A}$. Let
 269 us introduce the growth bound of $-\mathcal{A}$, given by

$$270 \quad (4.1) \quad \omega_0(-\mathcal{A}) := \lim_{t \rightarrow +\infty} t^{-1} \log (\|e^{-t\mathcal{A}}\|).$$

271 Assume that there exist $\mathcal{K} \in \mathcal{L}_c(\mathcal{H}, \mathcal{U})$, $C > 0$ and $\lambda > \omega_0(-\mathcal{A})$ such that for any
 272 $y^0 \in \mathcal{F}$, the solution to the abstract linear feedback system

$$273 \quad (4.2) \quad \begin{cases} y' &= \mathcal{A}y + \mathcal{BK}y, \\ y(0) &= y^0 \end{cases}$$

274 satisfies

$$275 \quad (4.3) \quad \forall t \geq 0, \|\mathcal{P}y(t)\|_{\mathcal{H}} \leq Ce^{-\lambda t}\|y^0\|_H.$$

276 Then, there exists $T > 0$ such that for any $y^0 \in \mathcal{F}$, there exists $u \in L^2((0, T), \mathcal{U})$ such
 277 that the solution y to the control system

$$278 \quad (4.4) \quad \begin{cases} y' &= \mathcal{A}y + \mathcal{B}u, \\ y(0) &= y^0 \end{cases}$$

279 satisfies $\mathcal{P}(y(T)) = 0$.

280 *Proof.* Our proof is inspired by [24, Theorem 2.4]. Assume that (4.3) holds. Let
 281 $T > 0$ and

$$282 \quad (4.5) \quad J(T) = \mathcal{P}e^{-T\mathcal{A}}e^{T(\mathcal{A}+\mathcal{BK})} \in \mathcal{L}_c(\mathcal{F}),$$

283 where $\mathcal{F}$ is endowed with the norm on $\mathcal{H}$. Moreover, since $\mathcal{P}$ commutes with $\mathcal{A}$,

$$284 \quad (4.6) \quad \|J(T)\| \leq \|e^{-T\mathcal{A}}\| \|\mathcal{P}e^{T(\mathcal{A}+\mathcal{BK})}\|.$$

285 (4.3) gives that as an operator from $\mathcal{F}$ to $\mathcal{F}$ endowed with the norm of $\mathcal{H}$, we have

$$286 \quad \|\mathcal{P}e^{T(\mathcal{A}+\mathcal{BK})}\| \leq Ce^{-\lambda T}.$$

287 Hence, the above inequality combined with (4.6) (recall that $\lambda > \omega_0(-\mathcal{A})$) ensures
 288 that $\|J(T)\| < 1$ for T large enough. Hence, setting $J = J(T)$ for such a fixed T , we
 289 have that $I - J$ is invertible and $(I - J)^{-1} \in \mathcal{L}_c(\mathcal{F})$.

290 Now, let $y^0 \in \mathcal{F}$. For $t \in [0, T]$, we set

$$291 \quad z_1(t) = e^{t(\mathcal{A}+\mathcal{BK})}(I - J)^{-1}y^0 \in H, \quad z_2(t) = e^{t\mathcal{A}}(I - (I - J)^{-1})y^0 \in H.$$

292 Seeing $(I - J)^{-1}y^0$ and $(I - (I - J)^{-1})y^0$ as elements of $\mathcal{H}$ and using a perturbation
 293 property of semigroups (see [28, Section 3.1, Formula (1.2)], we have

$$294 \quad (4.7) \quad y(t) := z_1(t) + z_2(t) = e^{t\mathcal{A}}y^0 + \int_0^t e^{(t-s)\mathcal{A}}\mathcal{BK}z_1(s)ds.$$

295 Clearly, $y(0) = y^0$. Moreover, from the expression of J given in (4.5), together with
 296 the fact that $\mathcal{P}$ and $\mathcal{A}$ commute, we deduce that

$$297 \quad \mathcal{P}e^{T(\mathcal{A}+\mathcal{BK})} = e^{T\mathcal{A}}J.$$

298 Using that $y^0 \in \mathcal{F} = \mathcal{P}(\mathcal{H})$, that $\mathcal{P}$ and $\mathcal{A}$ commute, and that $\mathcal{P}$ is a projection
 299 on $\mathcal{F}$, (so that $\mathcal{P}((I - J)^{-1}y^0) = ((I - J)^{-1}y^0)$), we obtain

$$\begin{aligned} \mathcal{P}(y(T)) &= \mathcal{P}\left(e^{T(\mathcal{A}+\mathcal{BK})}(I - J)^{-1}y^0\right) + \mathcal{P}\left(e^{T\mathcal{A}}(I - (I - J)^{-1})y^0\right) \\ &= e^{T\mathcal{A}}\left(J(I - J)^{-1} + (I - (I - J)^{-1})\right)y^0 \\ &= 0, \end{aligned}$$

whence the result by taking as a control

$$u(t) = \mathcal{K}z_1(t) \in L^2((0, T), \mathcal{U}), \quad \square$$

using (4.7) and the Duhamel formula.

Applying the previous proposition gives the following negative result.

PROPOSITION 4.4. (3.13) is not exponentially stable with respect to the semi-norm $\|\cdot\|_{H_{p,0}^{1/2}} \times \|\cdot\|_{L_{p,0}^2}$.

Proof. We argue by contradiction. We assume that for any $(\phi_{er}^0, \eta_{er}^0)$, the solution (ϕ_{er}, η_{er}) to (3.12) converges exponentially to 0 in the $\|\cdot\|_{H_{p,0}^{1/2}} \times \|\cdot\|_{L_{p,0}^2}$ semi-norm as $t \rightarrow +\infty$.

We can apply Lemma 4.3 with $\mathcal{A} = A$ (which is of growth bound 0), $\mathcal{B} = C$, $\mathcal{K} = L$ (note that L and C commute, so that $\mathcal{BK} = LC$) and $\mathcal{P}$ is the orthogonal projection on $\mathcal{F} = H_0 = H_{p,0}^{1/2} \times L_{p,0}^2$, which commutes with A (see the beginning of the proof of Theorem 4.1). We deduce that there exists $T > 0$ for which for any $(\phi^0, \eta^0) \in H_{p,0}^{1/2} \times L_{p,0}^2$, there exists $v \in L^2((0, T), \mathbb{T})$ such that the solution (ϕ, η) to

$$(4.8) \quad \begin{cases} \partial_t \phi(t, x) &= -g\eta(t, x), \\ \partial_t \eta(t, x) &= \mathcal{G}\phi(t, x) + \mathbf{1}_\omega(x)v(t, x), \\ \phi(0, x) = \phi^0(x), \quad \eta(0, x) &= \eta^0(x), \end{cases}$$

is such that $\mathcal{P}(\phi(T, \cdot), \eta(T, \cdot)) = 0$. Since $\mathcal{F}^\perp = E_0$ (the generalized eigenspace associated with 0), this implies that both $\phi(T, \cdot)$ and $\eta(T, \cdot)$ are constant.

Now, we introduce

$$u = \frac{i}{\sqrt{g}} \sqrt{\mathcal{G}}\phi + \eta.$$

Then, using (4.8), u solves

$$\begin{cases} \partial_t u &= -i\sqrt{g}\sqrt{\mathcal{G}}u + \mathbf{1}_\omega v, \\ u(0, x) &= \frac{i}{\sqrt{g}} \sqrt{\mathcal{G}}\phi^0(x) + \eta^0(x) \in L_{p,0}^2. \end{cases}$$

Clearly,

$$(\phi^0, \eta^0) \in H_{p,0}^{1/2} \times L_{p,0}^2 \mapsto \frac{i}{\sqrt{g}} \sqrt{\mathcal{G}}\phi^0(x) + \eta^0(x) \in L_{p,0}^2$$

is onto. We deduce that for any $u^0 \in L_{p,0}^2$, there exists

$$(\phi^0, \eta^0) \in H_{p,0}^{1/2} \times L_{p,0}^2$$

such that

$$u^0(x) = \frac{i}{\sqrt{g}} \sqrt{\mathcal{G}}\phi^0(x) + \eta^0(x).$$

For such ϕ^0, η^0 , there exists $v \in L^2((0, T) \times \mathbb{T})$ such that the solution (ϕ, η) to (4.8) is such that $\phi(T, \cdot)$ and $\eta(T, \cdot)$ are constants. Hence, posing u as above, we deduce that the solution u to

$$(4.9) \quad \begin{cases} \partial_t u &= -i\sqrt{g}\sqrt{\mathcal{G}}u + \mathbf{1}_\omega w, \\ u(0, x) &= u^0(x), \end{cases}$$

satisfies that $u(T, \cdot)$ is a constant.

Let us prove that we are in the setting of [21, Proposition 27], that enables to recover (under certain conditions) a full controllability result from a “partial” one, of the form set out above. Let us locally use the notations of this article, for the sake of clarity.

We set $H = L_p^2$, that is indeed a complex Hilbert space for the natural scalar product. We set $U = L_p^2$, $U_T = U = L^2((0, T), L_p^2)$. It satisfies the “extension by 0 property”: if $w \in U_T$, $a, b > 0$, then the function $\tilde{w}$ defined by $\tilde{w}(t) = 0$ for $0 < t < a$, $\tilde{w}(t) = w(t - a)$ for $a < t < T + a$, and $\tilde{w}(t) = 0$ for $T + a < t < T + a + b$ is in $L^2((0, T + a + b), L_p^2)$.

We also set $A = -i\sqrt{g}\sqrt{\mathcal{G}}$, which is an unbounded operator with domain $H_p^{\frac{1}{2}}$, $\mathcal{F} = \text{span}(1)$ (which is finite-dimensional and stable under e^{tA}), $\mathcal{S} = L_{p,0}^2 = \text{span}(1)^{\perp 1}$ (which is closed and finite codimensional), and $B = \mathbf{1}_\omega$. We have just proved that for any $u^0 \in \mathcal{S}$, there exists $w \in U_T$ such that the solution u of (4.9) satisfies that $u(T, \cdot) \in \mathcal{F}$.

Moreover, A is diagonalizable and its eigenfunctions are the family of the complex exponentials $\{e^{inx}\}_{n \in \mathbb{Z}}$, which is well-known to be a linearly independent family on any open subset of $\mathbb{T}$. Hence, since $B^* = \mathbf{1}_\omega$, we indeed have by [21, Remark 28] that for every finite linear combination of generalized eigenfunctions g_0 of A^* , we have $B^*e^{tA^*}g_0 = 0$ on $(0, \varepsilon)$ for any $\varepsilon > 0$ implies that $g_0 = 0$.

Hence, we can apply [21, Proposition 27] and obtain the null controllability of (4.8) on the whole state space L_p^2 at any time $T' > T$.

However, this turns out to be false. Indeed, using the formalism of [20], we can write

$$i\sqrt{g}\sqrt{\mathcal{G}} = \rho\left(\sqrt{-\Delta}\right),$$

with

$$\rho(x) = i\sqrt{g}\sqrt{x \tanh(x)}.$$

Note that ρ is holomorphic on $\mathbb{C}^*$ and continuous on $\mathbb{C}$, and that $\rho(z) = o(z)$ as $z \rightarrow \infty$. Then, we can apply [20, Theorem 1.4] to obtain that (4.9) is not controllable in L_p^2 , which yields the desired result by contradiction. $\square$

REMARK 4.5. *Contrary to Proposition 4.1 (and also to Proposition 5.1), we are not able to extend the previous result to more general observer gains $L \in \mathcal{L}_c(H)$ in the observation problem 3.15. Indeed, what is important in the previous reasoning is that we are able to see the output C as a control $\mathcal{B}$, and the observer gain L as a feedback $\mathcal{K}$, in order to apply Lemma 4.3, which makes a link between a partial stabilization result and a partial controllability result, but this lemma does not seem to be extendable to an observation problem. In fact, it is easy to figure out that Proposition 4.4 can be extended to any observer gain $L \in \mathcal{L}_c(H)$ which commutes with the output C , which is a highly restrictive condition.*

5. Modifying the state space and the Luenberger observer. As we have seen, we cannot have exponential stability of (3.13) because of two obstructions, one coming from the mean values and one coming from the high frequencies. Let us give two remedies to these problems.

¹We modified the notation $\mathcal{G}$ in [21, Proposition 27] into $\mathcal{S}$, to avoid confusions with the operator $\mathcal{G}$ of the present article.

5.1. State space of initial conditions with null mean value: polynomial decay. Let us restrict our state space to the set of initial conditions with null mean values, *i.e.* we take as a state space H_0 . Note that it is not very convenient to look at (3.12) with initial state in H_0 . Indeed, if we are looking at (3.13), we observe that H_0 is not stable under the semigroup generated by the operator appearing in (3.13): if we start with initial conditions in H_0 , then, in general, the term $\mathbf{1}_\omega \eta_{er}$ does not have mean value zero. One natural remedy is to “force” in some sense this property by projecting again on H_0 . We call Π_{H_0} the orthogonal projection in H onto the closed subspace H_0 . According to the previous discussion, we first propose a new Luenberger observer leading to the system

$$(5.1) \quad \begin{cases} \partial_t \hat{\phi}(t, x) &= -g \hat{\eta}(t, x), \\ \partial_t \hat{\eta}(t, x) &= \mathcal{G} \hat{\phi}(t, x) - \gamma \Pi_{H_0} \mathbf{1}_\omega(x) (\eta(t, x) - \hat{\eta}(t, x)), \\ \hat{\phi}(0, x) &= \hat{\phi}^0(x), \quad \hat{\eta}(0, x) = \hat{\eta}^0(x), \end{cases}$$

where $(\hat{\phi}^0(x), \hat{\eta}^0(x)) \in H_0$. If $(\phi^0(x), \eta^0(x)) \in H_0$ and (ϕ, η) is the corresponding solution to (3.4), then, $(\phi_{er}, \eta_{er}) := (\phi - \hat{\phi}, \eta - \hat{\eta})$ is the solution to

$$(5.2) \quad \begin{cases} \partial_t \phi_{er}(t, x) &= -g \eta_{er}(t, x), \\ \partial_t \eta_{er}(t, x) &= \mathcal{G} \phi_{er}(t, x) - \gamma \Pi_{H_0} (\mathbf{1}_\omega(x) \eta_{er}(t, x)), \\ \phi_{er}(0, x) &= \phi_{er}^0(x), \quad \eta_{er}(0, x) = \eta_{er}^0(x), \end{cases}$$

where $\phi_{er}^0(x) = \phi^0(x) - \hat{\phi}^0(x)$ and $\eta_{er}^0(x) = \eta^0(x) - \hat{\eta}^0(x)$ are in H_0 . For later purposes, we need to put our system under the form of a collocated control problem. An explicit computation shows that $\gamma \Pi_{H_0} \mathbf{1}_\omega$ is a self-adjoint and non-negative operator on H_0 . Hence, by functional calculus, it admits a unique square root denoted by $\tilde{B} = \sqrt{\gamma \Pi_{H_0} \mathbf{1}_\omega}$, which is also self-adjoint. Hence, this system can be rewritten in an abstract way as

$$(5.3) \quad \begin{cases} \partial_t \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix} &= (A - \tilde{B} \tilde{B}^*) \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix}, \\ \begin{pmatrix} \phi_{er}(0, x) \\ \eta_{er}(0, x) \end{pmatrix} &= \begin{pmatrix} \phi_{er}^0(x) \\ \eta_{er}^0(x) \end{pmatrix}. \end{cases}$$

Note that it can also be written in a different abstract way. Let us consider the output space $Y = H$ to be now different from the state space H_0 . Now, the output operator $\tilde{C}$ and the observer gain $\tilde{L}$ are given by

$$(5.4) \quad \tilde{C} = \begin{pmatrix} 0 & 0 \\ 0 & \mathbf{1}_\omega \end{pmatrix} \in \mathcal{L}_c(H_0, H) \text{ and } \tilde{L} = \gamma \Pi_{H_0} \in \mathcal{L}_c(H, H_0),$$

so that we also have

$$(5.5) \quad \begin{cases} \partial_t \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix} &= (A - \tilde{L} \tilde{C}) \begin{pmatrix} \phi_{er}(t, x) \\ \eta_{er}(t, x) \end{pmatrix}, \\ \begin{pmatrix} \phi_{er}(0, x) \\ \eta_{er}(0, x) \end{pmatrix} &= \begin{pmatrix} \phi_{er}^0(x) \\ \eta_{er}^0(x) \end{pmatrix}, \end{cases}$$

Moreover, remind that we already remarked that A is the generator of a C^0 -group on H_0 and that on H_0 , A is skew-adjoint. High frequencies are still an issue here, as shown in the next Proposition.

PROPOSITION 5.1. (5.3) is not exponentially stable with respect to the $\|\cdot\|_{H_{p,0}^{1/2}} \times \|\cdot\|_{L_{p,0}^2}$ -norm.

REMARK 5.2. Proposition 5.1 also holds if we replace $L = \gamma I_d$ by any $L \in \mathcal{L}_c(H)$, as is readily seen from the proof.

Proof. The proof is quite similar to that of Theorem 4.4. We argue by contradiction and we assume that (5.5) is exponentially stable. We deduce that the adjoint semigroup $(A - \tilde{L}\tilde{C})^* = -A - \tilde{C}^*\tilde{L}^*$ is exponentially stable. From the definition of $\tilde{C}$ given in (5.4), it is easy to deduce that $(\tilde{C})^* = \Pi_{H_0} \mathbf{1}_\omega \in \mathcal{L}_c(H, H_0)$. From [24, Theorem 2.4] (by seeing $(\tilde{C})^* = \Pi_{H_0} \mathbf{1}_\omega$ as a control operator and $(\tilde{L})^*$ as a feedback operator), we deduce that the system

$$\begin{cases} \partial_t \tilde{\phi}(t, x) &= g\tilde{\eta}(t, x), \\ \partial_t \tilde{\eta}(t, x) &= -\mathcal{G}\tilde{\phi}(t, x) + \gamma \Pi_{H_0}(\mathbf{1}_\omega(x)v(t, x)), \\ \tilde{\phi}(0, x) &= \phi^0(x), \quad \tilde{\eta}(0, x) = \tilde{\eta}^0(x), \end{cases}$$

is exactly controllable in the state space H_0 , at some time $T > 0$ for some $v \in L^2((0, T), L_p^2)$. Following the proof of Theorem 4.4, we argue by contradiction and we obtain that

$$\begin{cases} \partial_t u &= i\sqrt{g}\sqrt{\mathcal{G}}u + \Pi_{H_0}(\mathbf{1}_\omega v), \\ u(0, x) &= u^0(x) \end{cases}$$

is null controllable in the state space $L_{p,0}^2$ at some time $T > 0$, for some $v \in L^2((0, T), L_p^2)$. By duality (see e.g. [34, Theorem 11.2.1]), this means that there exists $C > 0$ such that for any $\varphi^0 \in L_{p,0}^2$, we have

$$\|\varphi(T, \cdot)\|_{L^2}^2 \leq \int_0^T \int_\omega |\varphi(t, x)|^2 dx dt,$$

where φ is the solution of the adjoint problem

$$\begin{cases} \partial_t \varphi &= -i\sqrt{g}\sqrt{\mathcal{G}}\varphi, \\ \varphi(0, x) &= \varphi^0(x). \end{cases}$$

It turns out that (5.7) is false. Indeed, the counterexample provided in [20, Proof of Theorem 1.4, Page 3145] is the periodization of a family of functions $(g_h)_{h \in (0,1)}$, defined on $\mathbb{R}$, that are compactly supported in the Fourier variable, with support away from $\xi = 0$. Notably, all the g_h are of mean 0 and are in L_p^2 , so they lie in $L_{p,0}^2$ and furnish a counterexample for (5.7). $\square$

However, even if one cannot expect exponential stabilization, we still have the following result.

THEOREM 5.1. (5.3) is asymptotically stable with respect to the $\|\cdot\|_{H_{p,0}^{1/2}} \times \|\cdot\|_{L_{p,0}^2}$ -norm, in the sense that the solution of (5.3) converges strongly to 0 as $t \rightarrow +\infty$. Moreover, (5.3) enjoys the following polynomial decay: there exists $C > 0$ such that for any $(\phi_{er}^0, \eta_{er}^0) \in D_0(A)$, the corresponding solution of (5.3) satisfies

$$\|(\phi_{er}(t, \cdot), \eta_{er}(t, \cdot))\|_{H_0} \leq \frac{C}{\sqrt{1+t}} \|(\phi_{er}^0, \eta_{er}^0)\|_{D_0(A)}, \quad t \geq 0.$$

REMARK 5.3. Remind that (5.8) cannot hold if we replace $\|(\phi_{er}^0, \eta_{er}^0)\|_{D_0(A)}$ by $\|(\phi_{er}^0, \eta_{er}^0)\|_{H_0}$ in the right-hand side, because of [13, Proposition V.1.7] (which asserts that in this case, (5.3) would be exponentially stable).

Proof. We would like to prove that [32, Theorem 3.6] applies. Here, $A^* = -A$ has compact resolvents, and the feedback is under the form $-\tilde{B}\tilde{B}^*$ (see (5.3)), which is exactly the form of the feedback constructed in [32, Theorem 3.6] (since it relies on [32, Theorems 3.3, 3.4, 3.5]). Moreover, what is called “complete controllability” in [32] is exactly what is called *approximate observability in infinite time* in [34, Definition 6.5.1]. Applying [34, Proposition 6.9.1], we deduce that approximate observability in infinite time holds as soon as the Fattorini-Hautus test holds: there does not exist any eigenfunction φ of the operator A^* such that $\tilde{B}^*\varphi = 0$, which is equivalent to $\tilde{B}\tilde{B}^*\varphi = 0$. We argue by contradiction. For such an eigenfunction $\varphi(x) = \begin{pmatrix} \phi(x) \\ \eta(x) \end{pmatrix}$ of $A^* = -A$ associated with an eigenvalue $\lambda \neq 0$, $\tilde{B}^*\varphi = 0$ is therefore equivalent to: for every $x \in \omega$, we have $\eta(x) = 0$. In view of the eigenfunctions given in (3.5), η is a linear combination of at most two distinct complex exponentials, that are well known to form a linearly independent family on any interval (and so on ω), we deduce that $\eta(x) = 0, \forall x \in \mathbb{T}$. In view of the expression of A , the fact that $A\varphi = \lambda\varphi$ notably gives (by looking at the first component) that $-g\eta = \lambda\phi$. Since $\lambda \neq 0$, we also obtain that $\phi(x) = 0, \forall x \in \mathbb{T}$, which concludes the proof of the asymptotic stabilization result by applying [32, Theorem 3.6].

Let us now explain how to prove (5.8). Our idea is to apply [5, Theorem 3.9]. In order to apply this Theorem, our first goal is to estimate

$$s\|\tilde{D}^*((1+is)^2 + \mathcal{G}_{-1})^{-1}\tilde{D}\|, s \in \mathbb{R}^+,$$

where $\tilde{D} = \tilde{D}^* = \sqrt{\gamma\Pi_{H_0}}1_\omega$, and $\mathcal{G}_{-1}$ is the extension of $\mathcal{G}$ in $\mathcal{L}_c(H_0, D_0(A)')$. Note that by [5, Section 2B], we have that for $s \geq 0$,

$$(s\tilde{D}^*((1+is)^2 + \mathcal{G}_{-1})^{-1}\tilde{D}) = \tilde{B}^*((1+is) - A_{-1})^{-1}\tilde{B},$$

where A_{-1} is the extension of A in $\mathcal{L}_c(L_{p,0}^2, (H_{p,0}^1)')$. Since $\tilde{B}^*$ is bounded, we have, for any $s \geq 0$,

$$\|\tilde{B}^*((1+is) - A_{-1})^{-1}\tilde{B}\| \leq \|\tilde{B}^*\| \|((1+is) - A_{-1})^{-1}\| \|\tilde{B}\|.$$

Since A has purely imaginary spectrum, we have the existence of $C > 0$ such that for any $s \geq 0$,

$$\|\tilde{B}^*((1+is) - A_{-1})^{-1}\tilde{B}\| \leq C.$$

Hence, coming back to (5.9), we have existence of some constant $C > 0$ such that for any $s \geq 0$, we have

$$s\|\tilde{D}^*((1+is)^2 + \mathcal{G}_{-1})^{-1}\tilde{D}\| \leq C.$$

Moreover, the eigenvalues $\sqrt{\mu_k}$ of $\sqrt{\mathcal{G}}$ behave asymptotically as $C\sqrt{k}$ as $k \rightarrow +\infty$ for some $C > 0$. In particular, there exists $c > 0$ small enough such that for any $k \in \mathbb{N}^*$, we have

$$\sqrt{\mu_{k+1}} - \sqrt{\mu_k} \geq \frac{c}{\sqrt{\mu_k}}.$$

For $s > 0$, we set

$$\delta_0(s) = \frac{c}{4s}.$$

Then, for any $s > 0$, the set $[s - \delta_0(s), s + \delta_0(s)]$ contains at most one eigenvalue. We call $WP(s)$ the “wavepacket” associated with $s \geq 0$, *i.e.* the spectral subspace associated with the set $[s - \delta_0(s), s + \delta_0(s)]$, which is either empty, or limited to an eigenspace of a unique eigenvalue.

Consider any eigenfunction f of $\sqrt{\mathcal{G}}$ associated with an eigenvalue $\sqrt{\mu_k}$. Then, f can be written as $f(x) = a_k e^{ikx} + b_k e^{-ikx}$ with $a_k, b_k \in \mathbb{C}$. Moreover, writing $\omega = (\alpha, \beta)$ with $\alpha < \beta$, using $\Pi_{H_0} f = f$, the fact that the scalar product on $L^2_{p,0}$ is the restriction of the scalar product on L^2_p , and (3.1),

$$\begin{aligned}
\|\tilde{D}^* f\|_{L^2_{p,0}}^2 &= \langle \tilde{D}^* f, \tilde{D}^* f \rangle_{L^2_{p,0}} \\
&= \langle \tilde{D} \tilde{D}^* f, f \rangle_{L^2_{p,0}} \\
&= \langle \Pi_{H_0} \mathbf{1}_\omega f, f \rangle_{L^2_p} \\
&= \langle \mathbf{1}_\omega f, \Pi_{H_0} f \rangle_{L^2_p} \\
&= \langle \mathbf{1}_\omega f, f \rangle_{L^2_p} \\
&= \frac{1}{2\pi} \int_\omega |f|^2 \\
&= \frac{1}{2\pi} \int_\alpha^\beta (|a|^2 + |b|^2 + 2 \operatorname{Re}(a\bar{b}e^{2ikx})) dx \\
&= \frac{\beta - \alpha}{2\pi} (|a|^2 + |b|^2) + \frac{1}{2\pi} \operatorname{Re} \left(a\bar{b} \int_\alpha^\beta e^{2ikx} dx \right) \\
&= \frac{\beta - \alpha}{2\pi} (|a|^2 + |b|^2) + \operatorname{Re} \left(a\bar{b} \frac{e^{2ik\beta} - e^{2ik\alpha}}{\pi ik} \right).
\end{aligned}$$

Since $\int_\mathbb{T} |f|^2 = |a|^2 + |b|^2$, we deduce that as soon as $f \neq 0$, we have $\|\tilde{D}^* f\|_{L^2_{p,0}}^2 \neq 0$ and, by Young’s inequality,

$$\frac{\|\tilde{D}^* f\|_{L^2_{p,0}}^2}{\int_\mathbb{T} |f|^2} \geq \frac{\beta - \alpha}{2\pi} - \frac{1}{\pi k} \rightarrow \frac{\beta - \alpha}{2\pi} > 0 \text{ as } |k| \rightarrow +\infty.$$

Hence, for some $C > 0$ small enough, we have that

$$(5.13) \quad \|\tilde{D}^* f\|_{L^2_{p,0}} \geq C \|f\|_{L^2_{p,0}}, \quad f \in WP(s), \quad s \geq 0.$$

Combining (5.11), (5.12) and (5.13), and applying [5, Theorem 3.9] gives that for some new $C > 0$

$$\|(isId - (A - \tilde{B}\tilde{B}^*))^{-1}\| \leq Cs^2, \quad s \in \mathbb{R}.$$

Applying [3, Theorem 2.4] gives that (for the operator norm in H_0)

$$(5.14) \quad \|e^{t(A - \tilde{B}\tilde{B}^*)}(A - \tilde{B}\tilde{B}^*)^{-1}\| = O\left(\frac{1}{\sqrt{t}}\right), \quad t \rightarrow +\infty.$$

We easily deduce (5.8) by applying this estimate together with the identity: $\forall \varphi \in D_0(A)$,

$$e^{t(A - \tilde{B}\tilde{B}^*)} \varphi = e^{t(A - \tilde{B}\tilde{B}^*)} (A - \tilde{B}\tilde{B}^*)^{-1} (A - \tilde{B}\tilde{B}^*) \varphi$$

and remarking that using a triangular inequality,

$$\|(A - \tilde{B}\tilde{B}^*) \varphi\| \leq C \|\varphi\|_{D_0(A)}. \quad \square$$

5.2. State space of low-frequency initial conditions with null mean value: exponential observer. A remedy is to look only at low-frequency functions. Let us fix some $N \in \mathbb{N}^*$. We call H_N the finite-dimensional space

$$H_N = \text{span}\{\nu_n^+, \nu_n^-\}_{1 \leq |n| \leq N}.$$

We also introduce $\Pi_{LF0^\perp}^N$ the orthogonal projection on H_N in H . We restrict ourselves to initial conditions in H_N and we propose the following observer:

(5.15)

$$\begin{cases} \partial_t \hat{\phi}(t, x) &= -g \hat{\eta}(t, x), \\ \partial_t \hat{\eta}(t, x) &= \mathcal{G}(\hat{\phi}(t, x)) - \gamma \Pi_{LF0^\perp}^N \mathbf{1}_\omega (\hat{\eta}(t, x) - \eta(t, x)), \quad (t, x) \in [0, +\infty) \times \mathbb{T}, \\ \hat{\phi}(0, x) &= \hat{\phi}^0(x), \quad \hat{\eta}(0, x) = \hat{\eta}^0(x), \end{cases}$$

resulting in the error equation:

$$(5.16) \quad \begin{cases} \partial_t \phi_{er}(t, x) &= -g \eta_{er}(t, x), \\ \partial_t \eta_{er}(t, x) &= \mathcal{G}(\phi_{er}(t, x)) - \gamma \Pi_{LF0^\perp}^N \mathbf{1}_\omega \eta_{er}(t, x), \quad (t, x) \in [0, +\infty) \times \mathbb{T}, \\ \phi_{er}(0, x) &= \phi_{er}^0(x), \quad \eta_{er}(0, x) = \eta_{er}^0(x), \end{cases}$$

where $\phi_{er}^0(x) = \phi^0(x) - \hat{\phi}^0(x)$ and $\eta_{er}^0(x) = \eta^0(x) - \hat{\eta}^0(x)$. We then have the following result.

THEOREM 5.2. (5.16) is exponentially stable, and there exists $C > 0$, independent of N , such that for any $(\phi_{er}^0(x), \eta_{er}^0(x)) \in H_N$ and any $t \geq 0$, we have

$$(5.17) \quad \|(\phi_{er}(t, \cdot), \eta_{er}(t, \cdot))\|_{H_0} \leq \|(\phi_{er}^0(x), \eta_{er}^0(x))\|_{H_0} e^{-\frac{C}{N}t}.$$

Proof. In this proof, $C > 0$ will be a constant, independent of N , that might change from line to line.

We introduce

$$B_N = \begin{pmatrix} 0 & 0 \\ 0 & \Pi_{LF0^\perp}^N \mathbf{1}_\omega \end{pmatrix},$$

which is a selfadjoint operator (since $\Pi_{LF0^\perp}^N \mathbf{1}_\omega = \Pi_{LF0^\perp}^N \mathbf{1}_\omega \Pi_{LF0^\perp}^N$ on H_N) and

$$\tilde{B}_N = \begin{pmatrix} 0 & 0 \\ 0 & \sqrt{\Pi_{LF0^\perp}^N \mathbf{1}_\omega} \end{pmatrix},$$

so that $B_N = \tilde{B}_N \tilde{B}_N^*$. Since A is skew-adjoint and B_N is selfadjoint on H_N , $A - B_N$ is dissipative in the sense that for any $\varphi \in H_N$, we have

$$\langle (A - B_N)\varphi, \varphi \rangle_H \leq 0.$$

Hence, we deduce that

$$(5.18) \quad \|e^{t(A-B_N)}\| \leq 1, \forall t \geq 0.$$

Following the proof of Theorem 5.1, we deduce by using a similar estimate to (5.10) that we have an estimate similar to (5.11), namely,

$$s \|\tilde{D}_N^* ((1 + is)^2 + \mathcal{G}_{-1})^{-1} \tilde{D}_N\| \leq C,$$

528 where $\tilde{D}_N = \sqrt{\Pi_{LF0^\perp}^N 1_\omega}$. Moreover, following exactly the notations and reasoning of
 529 the proof of [Theorem 5.1](#), we also have the analogous of [\(5.13\)](#)

$$530 \quad \|\tilde{D}_N f\|_{L_{p,0}^2} \geq C \|f\|_{L_{p,0}^2}, \quad f \in WP(s), \quad s \geq 0,$$

and the wavepacket is defined with the same δ_0 as in [\(5.12\)](#). We deduce that a result similar to [\(5.14\)](#), namely, for some $C_1 > 0$ large enough independant on N , we have

$$\|e^{t(A-B_N)}(A-B_N)^{-1}\varphi\|_H \leq \frac{C\|\varphi\|_H}{\sqrt{t}}, \quad \forall \varphi \in H_N, \quad \forall t \geq C_1,$$

531 which can be rewritten as

$$532 \quad (5.19) \quad \|e^{t(A-B_N)}\varphi\|_H \leq \frac{C\|(A-B_N)\varphi\|_H}{\sqrt{t}}, \quad \forall \varphi \in H_N, \quad \forall t \geq C_1.$$

533 By definition of H_N , since $\|B_N\| \leq C$, and the eigenvalues $\lambda_k^\pm$ of A behave asymp-
 534 totically like $\sqrt{|k|}$ as $|k| \rightarrow +\infty$, we have

$$\|(A-B_N)\varphi\|_H \leq C\sqrt{N}\|\varphi\|_H, \quad \forall \varphi \in H_N.$$

535 Hence, [\(5.19\)](#) implies that

$$536 \quad (5.20) \quad \|e^{t(A-B_N)}\varphi\|_H \leq \frac{C\sqrt{N}\|\varphi\|_H}{\sqrt{t}}, \quad \forall \varphi \in H_N, \quad \forall t \geq C_1.$$

Consider N large enough such that

$$\frac{C\sqrt{N}}{\sqrt{C_1}} \geq \frac{1}{2},$$

537 and consider $t_N \geq 0$ such that

$$538 \quad (5.21) \quad \frac{C\sqrt{N}}{\sqrt{t_N}} = \frac{1}{2}.$$

539 We have $t_N \geq C_1$, and from [\(5.20\)](#),

$$540 \quad (5.22) \quad \|e^{t_N(A-B_N)}\| \leq \frac{1}{2}.$$

541 Now, consider any $t \geq 0$, and write $t = at_N + b$, with $a \in \mathbb{N}$ and $0 \leq b < t_N$. Using
 542 [\(5.22\)](#) and [\(5.18\)](#), we obtain that

$$543 \quad (5.23) \quad \|e^{t(A-B_N)}\| \leq \|e^{t_N(A-B_N)}\|^a \|e^{b(A-B_N)}\| \leq \left(\frac{1}{2}\right)^a \leq \left(\frac{1}{2}\right)^{\frac{t}{t_N}}.$$

From [\(5.21\)](#), we deduce that $t_N = CN$. Hence, [\(5.23\)](#) gives that

$$\|e^{t(A-B_N)}\| \leq \left(\frac{1}{2}\right)^{\frac{C}{N}t},$$

544 from which we easily deduce [\(5.17\)](#). □

6. Numerical Results. This section is structured into four parts. The first part addresses certain numerical aspects that are worth discussing before proceeding with the numerical experiments. Specifically, it covers the distinction between performing computations in Fourier space versus classical space, the interpretation of low-frequency projectors, and the phenomenon of aliasing. The second part focuses on validating the theoretical results presented earlier through numerical simulations. The third part explores the behavior of the convergence rate as a function of the number of available frequencies. Finally, the fourth part shifts to an application beyond the theoretical framework of this paper. Here, we apply a Luenberger observer to the problem of wave field reconstruction in open water. Data assimilation for wave fields is a crucial topic in addressing certain challenges in engineering and hydrodynamics. This research area has gained significant attention in recent years, with applications ranging from flood or tsunami prediction [30, 26] and bathymetry detection [18, 2] to the reconstruction and prediction of wave fields from observations [35, 12, 27, 22, 8]. Our last case focuses on the latter, making simplified assumptions while ensuring that the proposed method remains both relevant and practical. The approach is deliberately straightforward and easy to implement.

6.1. Numerical aspects. In this section, we clarify some numerical aspects that are important to address in the following sections.

6.1.1. Fourier space and classical space. Let us first address the choice of the space (Fourier or classical) in which the numerical tests will be carried out. Computations for subsection 6.2 and subsection 6.3 are performed in Fourier space. Computations for subsection 6.4 are performed in the classical space.

Concerning subsection 6.2 and subsection 6.3, this choice is motivated by the desire to better represent the theoretical convergence results that we want to illustrate and to facilitate computations. But let us address the difference between running the simulations in the Fourier space and in the classical space. Recall that for a periodic continuous function f on the torus $\mathbb{T}$, there exists an equivalence between the representation of the function on the grid $(f(x_i))_{i=0,\dots,N_x-1}$ and its representation in the Fourier domain $\mathcal{F}[f](n)_{n=-N_f,\dots,N_f}$ if $N_f = (N_x - 1)/2$. This equivalence can be recovered with a simple discrete Fourier transform computation. However, this equivalence is no longer valid for piecewise continuous functions, which is our case because of the operator $\mathbf{1}_\omega$. In this setting we therefore lose the uniform and absolute convergence, but we still have convergence in L^2 -norm, and almost everywhere convergence of the Fourier series to the function in the classical space by classical theorems of Fourier analysis. It can be shown that the operators constructed subsequently in a discretized Fourier space converge towards the same operators constructed in a discretized classical space as N_x (or N_f) tends to infinity. The indicator operator $\mathbf{1}_\omega$ is then passed into Fourier space, assuming $\omega = [\pi - a, \pi + a]$, with:

$$(6.1) \quad c_0(\mathbf{1}_{[\pi-a, \pi+a]}) = \frac{1}{2\pi} \int_0^{2\pi} \mathbf{1}_{[\pi-a, \pi+a]}(x) dx = \frac{a}{\pi}$$

and, for $n \in \mathbb{Z} \setminus \{0\}$,

$$(6.2) \quad c_n(\mathbf{1}_{[\pi-a, \pi+a]}) = \frac{1}{2\pi} \int_0^{2\pi} \mathbf{1}_{[\pi-a, \pi+a]}(x) \exp(-inx) dx = \frac{a(-1)^n \sin(na)}{\pi na}.$$

Regarding subsection 6.4, computations are carried out in the classical space as we approach a realistic and practical problem.

6.1.2. Projectors. Let us clarify what *low-frequency projection* means in a numerical setting. Indeed, once we carry out a discretization (Fourier or classical), this already amounts to carrying out a low-frequency projection of our continuous problem into a finite-dimensional setting since a discretization of our function can only represent a finite number of frequencies.

This is what is considered in the numerical test performed in subsection 6.4. The *low-frequency projection* is not visible in the formulations, but rather implied by the fact that we work in a discretized setting.

Regarding the first numerical test in subsection 6.2, as our aim is there to support the theoretical assertions made in this paper, *low-frequency projection* is defined as the operator that projects the solution onto the frequency subspace that coincides with that of the initial condition. For example, if the initial condition of our initial value problem contains two sines of frequencies 1 and 2 (*i.e.*, $\sin(x) + \sin(2x)$), then the *low-frequency projection* operator projects a function in Fourier space onto frequencies $\{-2, -1, 0, 1, 2\}$. This initial condition is considered to be *low-frequency*.

6.1.3. Aliasing. Another important numerical aspect is aliasing [25]. It is also known as spectral folding and occurs when we try to represent a function that contains more frequencies than the discretization grid can allow. The frequencies that are too high for the grid are then folded onto the rest of the frequencies, hence introducing a numerical error that cannot be avoided.

In order to avoid aliasing in subsection 6.2, we choose to concentrate on an initial condition that can be represented on the chosen frequency space. These are therefore sums of sines and cosines, which contain fewer frequencies than the frequency space can represent.

In subsection 6.4, this phenomenon is present as we aim to apply our observer on a more realistic setting. We are then outside the theoretical frame exposed above.

6.2. Numerical convergence of the error equation in the theoretical setting. In this section, we validate numerically the theoretical results presented above. Computations are performed in MATLAB [17]. We choose to emulate a discretization in space of the 1D torus $[0, 2\pi[$ with N_x points, *i.e.* a $\{x_0, x_1, \dots, x_{N_x-1}\}$ grid with $x_i = i\Delta x$, where $\Delta x = 2\pi/N_x$. For simplicity, we assume N_x to be odd, so that the maximum number of frequencies that can be represented with this discretization is $N_f = (N_x - 1)/2$. We therefore suppose to use a grid with $N_x = 2^6 + 1$ equidistant points, equivalent to the same number of frequencies, ranging from $-N_f = (N_x - 1)/2$ to $N_f = (N_x - 1)/2$. Computations are actually carried out in this frequency space.

We want to compare the results of different ODEs that describe the behaviour of the error without modifications (6.3), the error with a projector that removes the mean (6.4), and the error with a projector that removes the mean and high frequencies (6.5).

$$(6.3) \quad [E_1] \Leftrightarrow \begin{cases} \partial_t \phi_{er}(t, x) = -g\eta_{er}(t, x) \\ \partial_t \eta_{er}(t, x) = \mathcal{G}(\phi_{er}(t, x)) - \mathbf{1}_\omega \eta_{er}(t, x) \end{cases}$$

$$(6.4) \quad [E_2] \Leftrightarrow \begin{cases} \partial_t \phi_{er}(t, x) = -g\eta_{er}(t, x) \\ \partial_t \eta_{er}(t, x) = \mathcal{G}(\phi_{er}(t, x)) - \Pi_{\perp 0} \mathbf{1}_\omega \eta_{er}(t, x), \end{cases}$$

$$(6.5) \quad [E_3] \Leftrightarrow \begin{cases} \partial_t \phi_{er}(t, x) = -g\eta_{er}(t, x) \\ \partial_t \eta_{er}(t, x) = \mathcal{G}(\phi_{er}(t, x)) - \Pi_{\perp 0} \Pi_{LF} \mathbf{1}_\omega \eta_{er}(t, x). \end{cases}$$

Here ω is defined as the interval $[\pi/2, 3\pi/2]$, representing an observation of half of the free surface.

Regarding the initial condition, recall that ϕ^0 and η^0 must satisfy relation (2.6), i.e., $\mathcal{F}[\eta^0](n) = \frac{i\omega_n}{g}\mathcal{F}[\phi^0](n)$ where $\mathcal{F}[\cdot]$ is the Fourier transform on $\mathbb{T}$. Same holds for $\hat{\eta}^0$ and $\hat{\phi}^0$ and consequently, for ϕ_{er}^0 and η_{er}^0 . We therefore choose to define ϕ_{er}^0 as a signal with no mean defined as a sum of $N_{init} = 4$ sines and cosines:

$$(6.6) \quad \phi_{er}^0(x) = \sum_{n=1}^{n=N_{init}=4} \beta_{n,1} \sin(nx) + \beta_{n,2} \cos(nx),$$

with $\beta_{n,1}$ and $\beta_{n,2}$ taken at random from $[0, 1]$ for each mode. Once ϕ_{er}^0 is set, we choose η_{er}^0 so that it follows relation (2.6) on each Fourier mode.

We discretize the operators associated with the dynamics of equations (6.3), (6.4) and (6.5) into matrices $\mathbf{M}_1$, $\mathbf{M}_2$ and $\mathbf{M}_3$ over $2N_f + 1$ frequencies, sorted as $\{0, 1, \dots, N_f, -N_f, \dots, -1\}$. This leads to:

$$(6.7) \quad \mathbf{M}_1 = \begin{pmatrix} \mathbf{O} & -g\mathbf{I} \\ \mathbf{G} & -\mathbf{C}_\omega \end{pmatrix}, \quad \mathbf{M}_2 = \begin{pmatrix} \mathbf{O} & -g\mathbf{I} \\ \mathbf{G} & -\mathbf{D}_{\perp 0}\mathbf{C}_\omega \end{pmatrix}, \quad \mathbf{M}_3 = \begin{pmatrix} \mathbf{O} & -g\mathbf{I} \\ \mathbf{G} & -\mathbf{D}_{LF}\mathbf{D}_{\perp 0}\mathbf{C}_\omega \end{pmatrix},$$

where $\mathbf{G} = \text{diag}(\{0, \tanh(d), \dots, N_f \tanh(dN_f), N_f \tanh(dN_f), \dots, \tanh(d)\})$ is the diagonal matrix corresponding to the Fourier multiplier operator $\mathcal{G}$, $\mathbf{I}$ the identity matrix, $\mathbf{O}$ the zero matrix, $\mathbf{C}_\omega$ the matrix representing in Fourier space the application of the convolution product of the Fourier coefficients of the operator $\mathbf{1}_\omega$ onto a vector, $\mathbf{D}_{\perp 0}$ the diagonal matrix which removes mode 0 and $\mathbf{D}_{LF}$ the *low-frequency* projection matrix onto frequencies $\{-N_{init}, \dots, N_{init}\}$. Note that, for simplicity reasons, we choose to set $\gamma = 1$. For a given matrix $\mathbf{M}_i$, $i \in \{1, 2, 3\}$, we solve the corresponding homogeneous linear ODE with *MATLAB* `ode45` solver. Numerical results are shown in Figure 6.1.

The non-convergence of the error equation (6.4) can therefore be observed, due to the obstruction coming from mode 0. Concerning the solution of the error equation (6.4), we can see the obstruction coming from the high frequencies, which prevents the desired convergence at low frequencies. Convergence is linear because, since the equation is discretized, we are in a finite-dimensional regime. But the rate of convergence is determined by the highest frequency represented by the grid, and not by the highest frequency which constitutes the initial condition. We can also see in Figure 6.1 (top right) and (bottom left) that although not existing in the initial conditions, high frequencies appear due to the indicator operator which mixes the frequencies, therefore supporting our theoretical results.

Finally, we can see that the solution to the error equation (6.5) converges at a linear rate determined by the highest frequency of the projector, which, as we know, is the same as in the initial condition. Figure 6.1 (bottom left) shows that this choice of observer does not create high frequencies.

6.3. Numerical study of the convergence rate. In this test, we focus on the rate of convergence of the observer. Once discretized in space, the problem falls within the framework of the section 5. We have linear convergence of the observer. This rate of convergence is therefore given by the largest real part among the eigenvalues of the observer matrix, i.e. the matrix $\mathbf{M}_2$ as presented in (6.7). We also know that this convergence rate depends on the highest frequency for a fixed observation interval ω . In Figure 6.2, we are therefore interested in the relationship between the number of

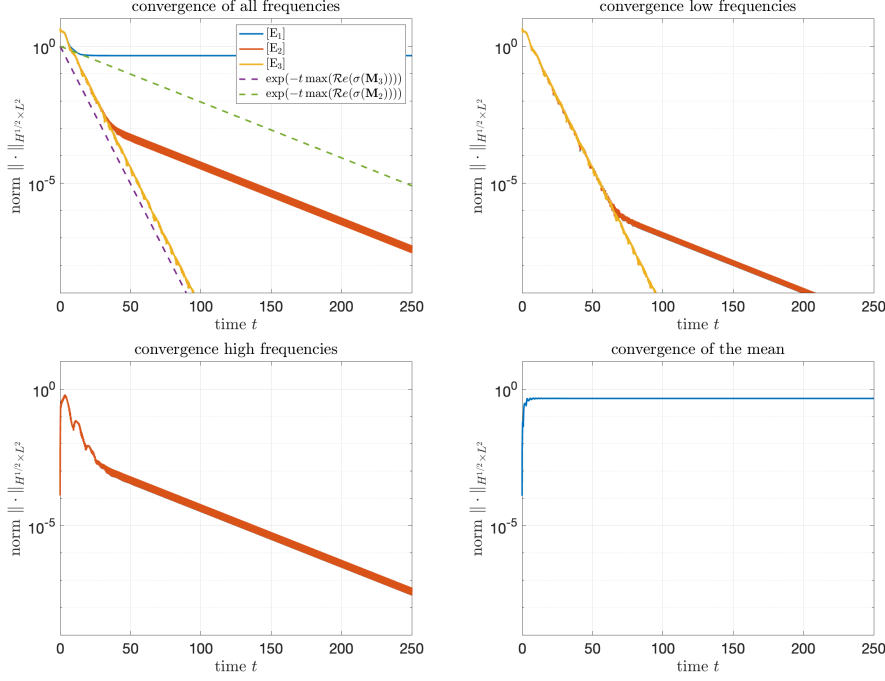


FIGURE 6.1. Convergence of the error equation for the initial condition in (6.6). Top left: Convergence of the solution of $[E_1]$, $[E_2]$ and $[E_3]$, with the theoretical rate of convergence found after an eigenvalue analysis of matrices $\mathbf{M}_2$ and $\mathbf{M}_3$. Top right: Same as top left, where the solutions have been projected onto the low frequencies, i.e., $\{-N_{init}, \dots, N_{init}\}$. Bottom left: Same as top left, where the solutions have been projected onto the high frequencies, i.e., $\{-N_f, \dots, N_f\} \setminus \{-N_{init}, \dots, N_{init}\}$. Bottom right: Same as top left, where the solutions have been projected onto the frequency 0.

676 frequencies represented and the rate of convergence. For increasing values of N_f , an
 677 increasingly large matrix $\mathbf{M}_2$ is constructed and its eigenvalue of largest real part is
 678 then extracted.

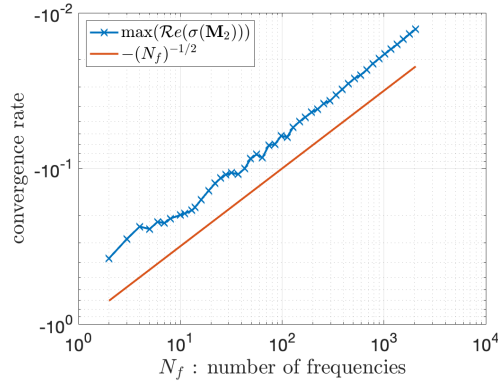


FIGURE 6.2. Convergence factor α as in $\|(\phi_{er}, \eta_{er})(t)\| \approx \exp(-\alpha t)$ as function of the number of frequencies N_f represented.

679 As can be seen in Figure 6.2, we therefore observe a linear convergence of the error

at rate $(N_f)^{-1/2}$. Recall that theoretically, by [Theorem 5.2](#), we have found that the convergence rate is at least of rate $(N_f)^{-1}$. Hence, it seems likely that our theoretical result is not optimal.

6.4. Application to wave field reconstruction. In this section, we focus on applying a Luenberger observer to a wave field in open sea, within a simplified but still relevant framework. These data are generated over a basin of length L .

$$(6.8) \quad \eta(t, x) = \operatorname{Re} \left\{ \sum_{n=1}^N A_n \exp(i(|k_n|x - \omega_n t + \lambda_n)) \right\}$$

where A_n are wave amplitudes of $N = 2048$ different individual waves calculated from a JONSWAP spectrum [\[16\]](#) with a peak period of 10 seconds, and a significant height of 3 meters. We then choose $L = 64l_p$ where l_p is the peak wavelength. We then have $k_n = \frac{n2\pi}{L}$, and $\omega_n = \sqrt{g|k_n| \tanh(d|k_n|)}$, and λ_n are random phase shifts drawn on $[0, 2\pi[$.

We then aim to set up an observer on the interval $\Omega = [0, L/4[$, an interval divided into two parts, with an observation region $\tilde{\Omega} = [0, L/8[$ and a prediction region $\Omega \setminus \tilde{\Omega} = [L/8, L/4[$ as illustrated in [Figure 6.3](#). We therefore simulate a situation where we observe a wave field over an interval of size $8l_p$, and where we reconstruct this wave field over an interval twice as large, with the aim of predicting it over a region where we have no information. Ω is then discretized with a grid of $N_\Omega = 256$ points (we therefore have $N_{\tilde{\Omega}} = N_{\Omega \setminus \tilde{\Omega}} = 128$). We then solve the following observer equations

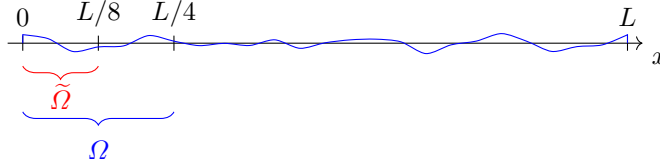


FIGURE 6.3. Setting for the reconstruction of a synthetic wave field.

on $\Omega \times [0, T_{\text{end}}]$:

$$(6.9) \quad \begin{cases} \partial_t \eta_o(x_l, t) &= -g\phi_o(x_l, t), \\ \partial_t \phi_o(x_l, t) &= \mathcal{G}(\eta_o(x_l, t)) + \gamma \Pi_{\perp 0} \mathbf{1}_{\tilde{\Omega}}(\eta(x_l, t) - \eta_o(x_l, t)), \\ \eta_o(x_l, 0) &= \mathbf{1}_{\tilde{\Omega}} \eta(x_l, 0), \\ \phi_o(x_l, 0) &= 0. \end{cases}$$

The theoretical results given in [section 5](#) no longer hold for two reasons. First, the η wave field is not periodic on Ω . Second is that the discretization of Ω is coarser than that used to generate η synthetically. As explained in [subsection 6.1](#), this second reason leads to aliasing, preventing us from achieving the linear convergence described in [section 5](#). However, as long as the spectrum support of η is largely covered by the discretization grid of Ω , we expect qualitative results. A prior statistical understanding of the spectrum of an open ocean wave field is therefore necessary, and this is why operations such as JONSWAP have been performed.

To validate our results, we use a different metric from the one used so far, called SSP (*Surface Similarity Parameter*), which is used in the hydrodynamics community

711 [29, 19, 9]. It is defined as follows, for two signals f_1 and f_2 :

$$712 \quad (6.10) \quad \text{SSP}(t) = \frac{(\sum_n |c_n(f_1) - c_n(f_2)|^2)^{1/2}}{(\sum_n |c_n(f_1)|^2)^{1/2} + (\sum_n |c_n(f_2)|^2)^{1/2}}.$$

713 The SSP has values in the range $[0, 1]$, where 0 corresponds to a complete match and
714 1 corresponds to a perfect mismatch between the two signals.

715 Solving equation (6.9) with $\gamma = 0.2$ gives the results shown in Figure 6.4. Here, for
716 a fixed time t , we compare the SSP score between the η function of the synthetic wave
717 field evaluated at the grid points in Ω and the solution of (6.9). The corresponding
snapshots are shown in Figure 6.5. The choice of $\gamma = 0.2$ is explained later.

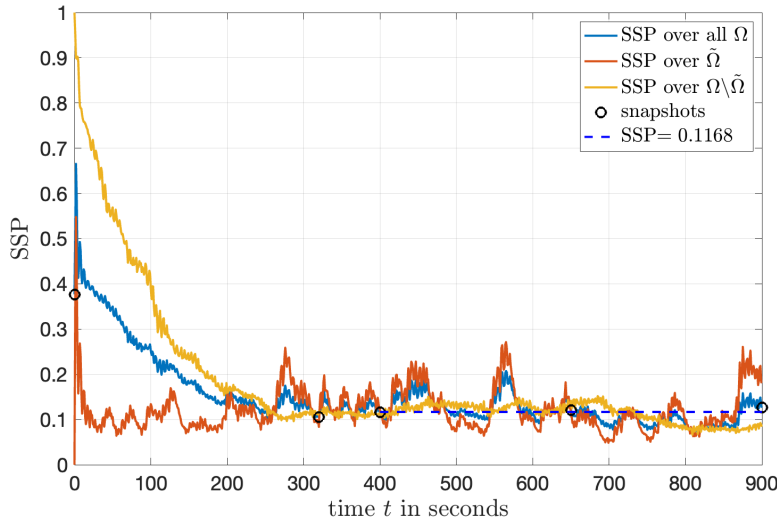


FIGURE 6.4. *SSP over time. Blue curve is the SSP on Ω , red curve is the SSP on $\tilde{\Omega}$, yellow curve is the SSP on $\Omega \setminus \tilde{\Omega}$, blue dotted curve is the mean value of the SSP on Ω over time interval $[400, 900]$.*

718 In Figure 6.4, we can see that the observer takes some time to reconstruct the
719 whole Ω interval (around 400 seconds of simulation) before converging around an
720 average SSP of 0.12. The snapshots in Figure 6.5 show a good reconstruction of the
721 surface over Ω from $t = 400$ seconds. We are then interested in the impact of the
722 gain constant γ on the observer, its convergence rate, and its average SSP once the
723 first 400 seconds of simulation have elapsed. The results are shown in Figure 6.6.
724 Figure 6.6 shows that, depending on the value chosen for the gain constant γ , the
725 average SSP varies, and varies differently if we look at the whole simulation interval,
726 only the observed one or only the reconstructed one. The results shown in Figure 6.4
727 and Figure 6.5 are therefore those with $\gamma = 0.2$ which equalizes the SSP over the
728 observation interval and the reconstruction interval, but we note that we can obtain
729 better results over one of these intervals individually if we choose another value of γ .
730

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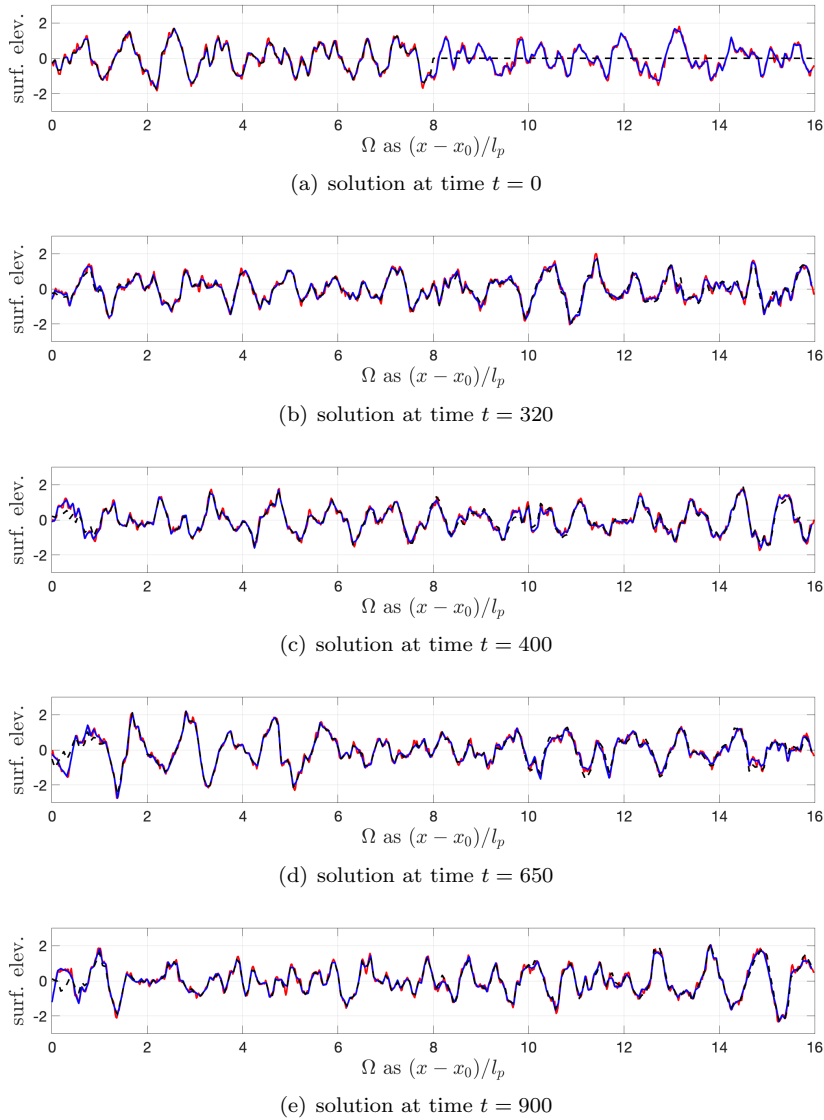


FIGURE 6.5. Snapshots of the surface over Ω at different timesteps. Red curve is the full synthetic solution η on a fine grid, blue curve is the synthetic solution η on the grid used for the observer, black dotted curve is the reconstructed surface η .

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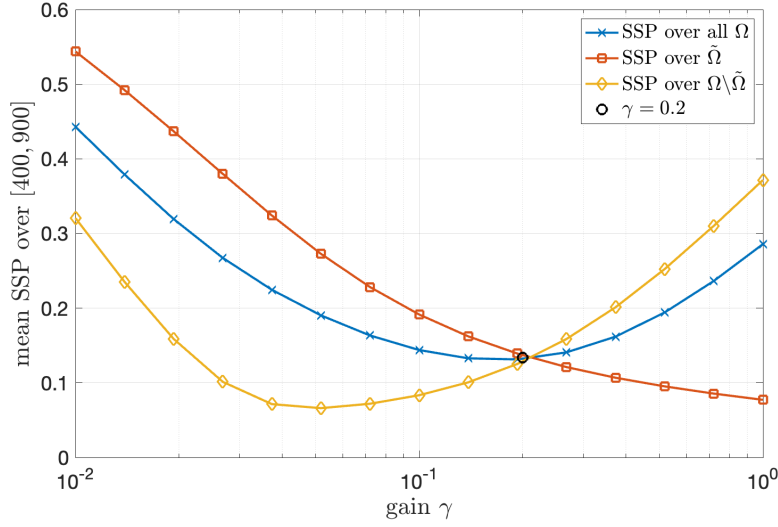


FIGURE 6.6. mean SSP over [400, 900] as a function of gain γ . Blue curve is the SSP on Ω , red curve is the SSP on $\hat{\Omega}$, yellow curve is the SSP on $\Omega \setminus \hat{\Omega}$.

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