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Continuous data assimilation for the 3D Ladyzhenskaya model: analysis and computations



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ABSTRACT

We analyze continuous data assimilation by nudging for the 3D Ladyzhenskaya equations. The analysis provides conditions on the spatial resolution of the observed data that guarantee synchronization to the reference solution associated with the observed, spatially coarse data. This synchronization holds even though it is not known whether the reference solution, with initial data in L^2 , is unique; any particular reference solution is determined by the observed, coarse data. The efficacy of the algorithm in both 2D and 3D is demonstrated by numerical computations.

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1. Introduction

1.1. Data assimilation

The insertion of coarse grain observational measurements into a mathematical model is called continuous data assimilation. This can provide a more accurate forecast in applications ranging from the medical, environmental and biological sciences, [1,2], to imaging, traffic control, finance and oil exploration [3]. Bayesian and variational approaches (Kalman filters, 3DVar and 4DVar) are based on discrete observations in time and often used to treat errors in both observed data and model itself [4–11]. They are widely used in practice, but difficult to analyze mathematically, especially for physical models governed by nonlinear differential equations [12–14].

Nudging is a straightforward, deterministic approach to data assimilation. While its origin can be traced back to [15], it has been more recently applied in the context of synchronizing chaotic dynamical systems.

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See [16,17] for a more complete history, and [18] for a comparison with Kalman filtering. In essence, this method assumes that an accurate initial condition $\mathbf{u}(0)$ is not known for a particular model

$$\frac{\mathrm{d}\mathbf{u}}{\mathrm{d}t} = \mathbf{F}(\mathbf{u}),\tag{1.1}$$

but data from a reference solution, interpolated at spatial resolution h is available, denoted as $I_h \mathbf{u}(t)$. Those observations are used in an auxiliary system

$$\frac{\mathrm{d}\mathbf{v}}{\mathrm{d}t} = \mathbf{F}(\mathbf{v}) - \mu I_h(\mathbf{v} - \mathbf{u}) , \quad \mathbf{v}(0) = 0$$
(1.2)

to drive $\|\mathbf{v} - \mathbf{u}\| \to 0$ at an exponential rate, provided μ is sufficiently large and h is sufficiently small. Since derivatives are not required of the data in the approach, it can be used with more common types of observations, such as nodal values.

Rigorous analysis of the nudging algorithm for partial differential equations in fluid mechanics began with the work of Azouani, Olson and Titi. They estimated threshold values for the relaxation parameter μ and data resolution h for the 2D NSE [19]. Since then, nudging has been rigorously shown to synchronize with reference solutions in a variety of applications, including the 2D Rayleigh–Bénard problem [20–22], surface quasigeostrophic equation [23,24], Korteweg–de Vries equation [25], 2D magnetohydrodynamic system [26], 3D Brinkman-Forchheimer-extended Darcy model [27], 3D primitive equations [28], 3D Leray- α model [29], and Voigt-relaxation of the 2D NSE [30].

In each case, threshold values for μ and h had to be established for both the well-posedness of the corresponding system (1.2) as well as for synchronization. In some works it has been shown that it is sufficient to nudge with data in only a subset of the system variables [20–22,31,32]. While the nudging algorithm does not lend itself to directly treat error in the model, the effect of error in the observed data has been studied in [24,33].

1.2. The Ladyzhenskaya model

The motion of an homogeneous, incompressible, viscous fluid in a domain $\Omega \subset \mathbb{R}^3$ is classically described by the momentum equation and the incompressibility constraint, that read as

$$\partial_t \mathbf{u} + (\mathbf{u} \cdot \nabla) \mathbf{u} - \nabla \cdot \mathbf{T}(\mathbf{D}\mathbf{u}) + \nabla P = \mathbf{f},$$

$$\nabla \cdot \mathbf{u} = 0,$$
in $\Omega \times (0, \infty),$
(1.3)

where \mathbf{u} is the fluid velocity, P is the fluid pressure, \mathbf{f} is the forcing term. Here, $\mathbf{D}\mathbf{u}$ denotes the symmetric part of the gradient of \mathbf{u} . In particular, the Navier–Stokes model corresponds to the case of Newtonian fluids characterized by the (linear) Stokes' law $\mathbf{T}(\mathbf{D}\mathbf{u}) = \nu_0 \mathbf{D}\mathbf{u}$. The lack of a global regularity result makes the analysis of the nudging algorithm problematic for the 3D NSE, though a recent work provides a condition on observed data which deals with this issue [34]. In this work we consider a family of 3D globally well-posed modified Navier–Stokes equations, namely the Ladyzhenskaya model. In the mid-1960s, a number of modifications to the Navier–Stokes equations were suggested by Ladyzhenskaya for the description of the dynamics of viscous fluids when velocity gradients are large [35–37]. These equations form an important mathematical model describing the flow behavior of a wide class of non-Newtonian fluids [38–40]. In this work we consider one particular model (see Eq. (3.1)), where the Cauchy tensor in (1.3) takes the following nonlinear form

$$\mathbf{T}(\mathbf{D}\mathbf{u}) = \left(2\nu_0 + 2\nu_1 |\mathbf{D}\mathbf{u}|_F^{p-2}\right) \mathbf{D}\mathbf{u},\tag{1.4}$$

where $D\mathbf{u} = \frac{1}{2} \left[(\nabla \mathbf{u}) + (\nabla \mathbf{u})^T \right]$ and the Frobenius norm $|D\mathbf{u}|_F = (D\mathbf{u} : D\mathbf{u})^{\frac{1}{2}}$. The above relation is commonly used for non-Newtonian fluids with shear dependent viscosity, i.e. the dynamic viscosity depends

on $|\mathrm{D}\mathbf{u}|_F^2$. The model corresponding to p=2 reduces to the Navier–Stokes equations (NSE) with kinematic viscosity equal to $(\nu_0 + \nu_1)$. For p=3, it is mathematically equivalent to the Smagorinsky model [41] and the NSE with the von Neumann Richtmyer artificial viscosity for shocks [42].

There are various reasons to consider the Ladyzhenskaya models instead of the Navier–Stokes equations. In the first place, the laws of conservation of mass and momentum provide an undetermined system of partial differential equations for the velocity, pressure, and stress tensor. In general, this system of equations is not closed until the stress tensor, which represents all the internal forces, is related to the fluid velocity. Internal forces, and therefore also the stress tensor, must depend on local velocity differences and some combination of derivatives of velocity, i.e., the deformation tensor. The simplest relation is a linear law between stress and deformation, which leads to the Navier–Stokes equations, see [43] for details. This linear relation is only an approximation for a real fluid and schematically the stress and deformation are related nonlinearly, especially for large deformations. A specific nonlinear mathematical relationship between the stress and deformation can be derived from Stokes' hypotheses. If one retains some of these nonlinear terms, we arrive at the Ladyzhenskaya model considered here. We refer the interested reader to [38–40] for more details.

Secondly and more from a practical engineering point of view, the study of the Ladyzhenskaya equations is related to the field of turbulence modeling. For some values of p, such as p=3,4, the Ladyzhenskaya model considered here is equivalent to those of popular turbulence models, such as large eddy simulation (LES) and zero-equation models. In both applications and turbulence modeling, the behavior of averaged quantities are most important and often simulated. To do so, the quantities describing the flow are decomposed into its averaged and fluctuating quantities. However, averaging the NSE yields a non-closed system; to close the system, one must provide the relationship between the fluctuating and the averaged quantities. There are a wide range of closure assumptions which are known collectively as turbulence closure models. Two examples that are widely used are the zero equation model (or algebraic model) and large eddy simulation, for more details see e.g., [44,45]. The main feature of these models is that the non-closed part (known as the Reynolds stresses), which represents the contribution of small scales in the system, is related to the derivatives of the averaged quantities. See [46,47] for more details on the mathematics of large eddy simulation.

Finally, from the theoretical point of view, while the well-posedness has not been proven for the Navier–Stokes equations in three space dimensions, several results of existence, uniqueness and regularity of global-in-time solutions of the Ladyzhenskaya model have been proved in the last decades [35–38,40,48–53]. This provides a firm mathematical foundation for the study of (3.1).

1.3. Results in this paper

In this paper we develop a comprehensive study based on the theoretical analysis and large eddy simulation of the nudged system (1.2) corresponding to the Ladyzhenskaya model (1.3)–(1.4) (see (4.1)) with both no-slip and periodic boundary conditions.

In the no-slip case, we first use the Schauder fixed point theorem to prove that the nudged system has a unique global weak solution with $\mathbf{u}_0 \in L^2_{\sigma}(\Omega)$ provided $p \geq 5/2$ (see Theorem 4.1). To prove the existence of the solution, unlike some treatments of the nudged system for other models (e.g. [19,54]), this approach does not require μ to be large, nor h to be small. Then, we find a threshold value μ^* in terms of p, ν_0 , ν_1 , domain size and the Grashof number (see (3.5)), such that for $\mu \geq \mu^*$ and h correspondingly small, synchronization is guaranteed. That is, the nudged solution \mathbf{v} converges exponentially fast to the reference solution \mathbf{u} after a transient (see Theorem 4.2). In the case of periodic boundary conditions, the existence of global weak solutions to the Ladyzhenskaya model over the wider range $p \geq 11/5$, originating from $\mathbf{u}_0 \in L^2_{\sigma}(\Omega)$, has been established in [38]. These weak solutions are not known to be unique, unless the initial condition is

¹ Stokes introduced a series of requirements which together serves to define a ordinary fluid such as water and air [43].

more regular (see [38, Theorem 4.37] for $\mathbf{u}_0 \in H^1_{\sigma}(\Omega)$, and also [48] for $\mathbf{u}_0 \in W^{1,p}_{\sigma}(\Omega)$ in the no-slip case). Nonetheless, any one such weak solution becomes more regular after some time. Following the strategy devised in [38,50], we prove, for the endpoint case p=11/5, a time averaged bound in $L^{\frac{11}{5}}(\bar{t},\infty;W^{1,\frac{33}{5}}(\Omega))$ for the solution \mathbf{u} to the Ladyzhenskaya model in terms of the parameters of the system, where \bar{t} is suitably chosen (see Theorem 5.1). This bound is then used to prove the synchronization of the nudged solution to the solutions of the Ladyzhenskaya model (see Theorem 5.3). In contrast to the previous cases studied in literature when uniqueness of the reference solution holds, the novelty of our result is that synchronization takes place even without uniqueness of the reference solution for $11/5 \leq p \leq 5/2$. More precisely, for each reference solution \mathbf{u} corresponding to an initial datum $\mathbf{u}_0 \in L^2_{\sigma}(\Omega)$, the nudged solution $\mathbf{v}_{\mathbf{u}}$ converges to \mathbf{u} at an exponential rate for large time. As a consequence of our analysis, it is worth concluding that if two reference solutions \mathbf{u} and $\tilde{\mathbf{u}}$ are such that $I_h(\mathbf{u})(t) = I_h(\tilde{\mathbf{u}})(t)$ as $t \to \infty$, then $\|\mathbf{u} - \tilde{\mathbf{u}}\| \to 0$ as $t \to \infty$, i.e., the model has a finite number of determining modes for $p \geq 11/5$.

We demonstrate the efficacy of the algorithm by extensive computational studies. Numerical work with other fluid systems has shown that the nudging algorithm achieves synchronization with data that is much more coarse than required by the rigorous estimates [55–58]. We find this is also the case for the Ladyzhenskaya model with periodic boundary conditions, for which we achieve exponential convergence to machine precision with $h \approx 0.1$. Most of our computations are done for the case where p=3 (Smagorinsky model), corresponding to large eddy simulation [41]. Though for periodic boundary conditions we present the analysis for a threshold value of μ only in the endpoint case p=11/5, our numerical computations, show virtually no sensitivity to p for two choices of μ . Finally we test an abridged nudging scheme which uses data only for the horizontal components of velocity. We present evidence that synchronization still holds for that scheme, though at a slower rate for the third component of velocity and pressure.

Organization of this paper

In Section 2, we introduce the inequalities and preliminary results used in the analysis. Section 3 provides background on the Ladyzhenskaya model. Later, in Sections 4 and 5, we state and prove our main results, in which we give conditions under which the approximate solutions, obtained by the data assimilation algorithm, converge to the solution of the Ladyzhenskaya equations. Numerical experiments, demonstrating and extending beyond the analytical results, are described in Section 6.

2. Notation and preliminaries

Let $\Omega \subset \mathbb{R}^d$, d = 2, 3, be a bounded open Lipschitz domain with volume $|\Omega|$ and let $p \in [1, \infty]$. The Lebesgue space $L^p(\Omega)$ is the space of all measurable functions \mathbf{v} on Ω for which

$$\|\mathbf{v}\|_{L^p} := \left(\int_{\Omega} |\mathbf{v}(\mathbf{x})|^p d\mathbf{x}\right)^{\frac{1}{p}} < \infty \quad \text{if } p \in [1, \infty),$$
$$\|\mathbf{v}\|_{L^{\infty}} := \underset{\mathbf{x} \in \Omega}{\operatorname{ess sup}} |\mathbf{v}(\mathbf{x})| < \infty \quad \text{if } p = \infty.$$

The L^2 norm and inner product will be denoted by $\|\cdot\|$ and (\cdot, \cdot) , respectively. Let **V** be a Banach space of functions defined on Ω with the associated norm $\|\cdot\|_{\mathbf{V}}$. We denote by $L^p(a,b;\mathbf{V})$, the Bochner space of measurable functions $\mathbf{v}:(a,b)\to\mathbf{V}$ such that

$$\begin{aligned} \|\mathbf{v}\|_{L^p(a,b;\mathbf{V})} &\coloneqq \left(\int_a^b \|\mathbf{v}(t)\|_{\mathbf{V}}^p \, dt\right)^{\frac{1}{p}} < \infty & \text{if } p \in [1,\infty), \\ \|\mathbf{v}\|_{L^\infty(a,b;\mathbf{V})} &\coloneqq \mathop{\mathrm{ess\,sup}}_{t \in (a,b)} \|\mathbf{v}(t)\|_{\mathbf{V}} < \infty & \text{if } p = \infty. \end{aligned}$$

The space $W_0^{1,p}(\Omega)$ consists of all functions in $W^{1,p}(\Omega)$ that vanish on the boundary $\partial\Omega$ (in the sense of traces)

$$W_0^{1,p}(\Omega) = \{ \mathbf{v} : \mathbf{v} \in W^{1,p}(\Omega) \text{ and } \mathbf{v}|_{\partial\Omega} = 0 \}.$$

We introduce the Banach spaces of solenoidal functions

$$L^2_{\sigma}(\Omega) = \{ \mathbf{v} : \mathbf{v} \in L^2(\Omega), \quad \nabla \cdot \mathbf{v} = 0 \text{ and } \mathbf{v} \cdot \mathbf{n}|_{\partial \Omega} = 0 \},$$

$$W^{1,p}_{\sigma}(\Omega) = \{ \mathbf{v} : \mathbf{v} \in W^{1,p}(\Omega), \quad \nabla \cdot \mathbf{v} = 0 \text{ and } \mathbf{v}|_{\partial \Omega} = 0 \},$$

which are equipped with the same norms as $L^2(\Omega)$ and $W_0^{1,p}(\Omega)$, respectively. The spaces $\dot{L}^2(\Omega)$, $\dot{L}^2_{\sigma}(\Omega)$, $\dot{H}^1_{\sigma}(\Omega)$, $\dot{W}^{1,p}_{\sigma}(\Omega)$ will consist of the subsets of $L^2(\Omega)$, $L^2_{\sigma}(\Omega)$, $H^1_{\sigma}(\Omega)$ and $H^1_{\sigma}(\Omega)$, respectively, whose functions have zero spatial average, i.e. $\int_{\Omega} \mathbf{u} \, d\mathbf{x} = 0$. We denote by $(W_{\sigma}^{1,p}(\Omega))'$ the dual space of $W_{\sigma}^{1,p}(\Omega)$. We recall the following inclusions for $p \geq 2$

$$W^{1,p}_\sigma(\varOmega)\subset L^2_\sigma(\varOmega)\subset \left(W^{1,p}_\sigma(\varOmega)\right)'\quad \text{if}\quad \frac{2d}{d+2}\leq p<\infty,$$

where these injections are continuous, dense and compact. For matrix $A = (a_{ij})_{i,j=1}^3$, the Frobenius norm of the matrix A is given by

$$|A|_F = \left(\sum_{i,j=1}^3 (a_{ij})^2\right)^{\frac{1}{2}} = (A:A)^{\frac{1}{2}}.$$

The data assimilation method requires that the observational measurements $I_h(u)$, with h > 0, be given as linear interpolant observables satisfying $I_h : L^2(\Omega) \to L^2(\Omega)$ such that

$$||I_h \varphi|| \le c_I ||\varphi||, \qquad \forall \varphi \in L^2(\Omega),$$

$$||\varphi - I_h \varphi|| \le c_0 h ||\varphi||_{H^1(\Omega)}, \qquad \forall \varphi \in H^1(\Omega).$$
 (2.1)

One example of such interpolation operators includes projection onto Fourier modes with wave numbers $|k| \le 1/h$. Somewhat more physical are the volume elements and constant finite element interpolation [59,60].

Inequalities in Banach and Hilbert spaces

We recall here some well-known inequalities in Banach and Hilbert spaces which can be found in the classical literature (see, e.g., [61,62]). Let $1 \le p \le \infty$, we denote by p' the conjugate exponent, $\frac{1}{p} + \frac{1}{p'} = 1$. Assume that $f \in L^p$ and $g \in L^{p'}$ with $1 \le p \le \infty$. Then

$$||fg||_{L^1} \le ||f||_{L^p} ||g||_{L^{p'}}.$$
 (Hölder inequality)

Moreover, for any $a, b \ge 0$ and $\lambda > 0$ we have

$$ab \le \lambda a^p + (p\lambda)^{-\frac{p'}{p}} \frac{1}{p'} b^{p'}.$$
 (Young inequality)

Suppose $1 , there exist two constants <math>c_P$ and c_K such that for any $\mathbf{f} \in W_0^{1,p}(\Omega)$

$$\|\mathbf{f}\|_{L^p} \le c_{\mathcal{P}} \|\nabla \mathbf{f}\|_{L^p},$$
 (Poincaré inequality)

and

$$\|\nabla \mathbf{f}\| \le \sqrt{2} \|\mathbf{D}\mathbf{f}\| \quad \text{if } p = 2, \qquad \|\nabla \mathbf{f}\|_{L^p} \le c_{\mathbf{K}} \|\mathbf{D}\mathbf{f}\|_{L^p} \quad \text{if } p \ne 2,$$
 (Korn inequality)

where $\mathrm{D}\mathbf{f} = \frac{1}{2}\left[\left(\nabla\mathbf{f}\right) + \left(\nabla\mathbf{f}\right)^T\right]$. The constants c_P and c_K depend only on p and Ω . In the sequel, we will make use of the classical embedding theorems for Sobolev spaces

$$W^{1,2}(\Omega) \hookrightarrow L^6(\Omega),$$
 (Sobolev embedding) $W^{1,3}(\Omega) \hookrightarrow L^p(\Omega), \quad \forall \, p \in [1, \infty).$

We recall the interpolation inequalities for Lebesgue and Sobolev spaces. Let $\mathbf{f} \in L^p \cap L^q$, with $1 \leq p, q \leq \infty$. Then, for all r such that

$$\frac{1}{r} = \frac{\theta}{p} + \frac{1-\theta}{q}, \quad 0 \le \theta \le 1,$$

it follows that $\mathbf{f} \in L^r$ and

$$\|\mathbf{f}\|_{L^r} \le \|\mathbf{f}\|_{L^p}^{\theta} \|\mathbf{f}\|_{L^q}^{1-\theta}$$
. (Lebesgue interpolation inequality)

In addition, for any $\mathbf{f} \in W_0^{1,2}(\Omega)$, we have

$$\|\mathbf{f}\|_{L^4} \le C_L \|\mathbf{f}\|^{\frac{1}{4}} \|\nabla \mathbf{f}\|^{\frac{3}{4}}.$$
 (Ladyzhenskaya's inequality)

We need the following fundamental results (see, e.g., [61,63]).

Theorem 2.1 (Schauder Fixed-point Theorem). Let \mathcal{X} be a Banach space, and let \mathcal{A} be a nonempty closed convex set in \mathcal{X} . Let $\mathcal{F}: \mathcal{A} \to \mathcal{A}$ be a continuous map such that $\mathcal{F}(\mathcal{A}) \subset \mathcal{K}$, where \mathcal{K} is a compact subset of \mathcal{A} . Then \mathcal{F} has a fixed point in \mathcal{K} .

Theorem 2.2 (Aubin–Lions–Simon). Let $\mathcal{B}_0 \subset \mathcal{B}_1 \subset \mathcal{B}_2$ be three Banach spaces. We assume that the embedding of \mathcal{B}_1 in \mathcal{B}_2 is continuous and that the embedding of \mathcal{B}_0 in \mathcal{B}_1 is compact. For $1 \leq p, r \leq +\infty$ and T > 0, we define

$$\mathbf{E}_{p,r} = \{ v \in L^p(0,T; \mathcal{B}_0) , v_t \in L^r(0,T; \mathcal{B}_2) \}.$$

Then, we have

- (1) If $p < +\infty$, the embedding of $\mathbb{E}_{p,r}$ in $L^p(0,T;\mathcal{B}_1)$ is compact.
- (2) If $p = +\infty$ and r > 1, the embedding of $\mathbb{E}_{p,r}$ in $\mathcal{C}(0,T;\mathcal{B}_1)$ is compact.

Lastly, we report the following Gronwall lemmas which will play a crucial role in our analysis (see, e.g., [64]).

Lemma 2.3 (Gronwall's Lemma in Differential Form). Let $T \in \mathbb{R}^+$, $f \in W^{1,1}(0,T)$ and $g, \lambda \in L^1(0,T)$. Then

$$f'(t) \le \lambda(t) f(t) + g(t)$$
 a.e. in $[0, T]$

implies for almost all $t \in [0, T]$

$$f(t) \le f(0) e^{\int_0^t \lambda(\tau) d\tau} + \int_0^t g(s) e^{\int_s^t \lambda(\tau) d\tau} ds.$$

Lemma 2.4 (Uniform Gronwall Lemma - 1). Let $T \in \mathbb{R}^+$, $f \in W^{1,1}(t_0, \infty)$ and $g, \lambda \in L^1_{loc}(t_0, \infty)$ which satisfy

$$f'(t) \le \lambda(t) f(t) + g(t)$$
 a.e. $in (t_0, \infty)$,

and

$$\int_t^{t+r} \lambda(\tau) d\tau \le a_1, \quad \int_t^{t+r} g(\tau) d\tau \le a_2, \quad \int_t^{t+r} f(\tau) d\tau \le a_3, \qquad \forall t \ge t_0,$$

for r, a_1 , a_2 and a_3 positive. Then, for r > 0, we have

$$f(t) \le \left(\frac{a_3}{r} + a_2\right) e^{a_1}, \qquad \forall t \ge t_0 + r.$$

Lemma 2.5 (Uniform Gronwall Lemma - 2). Let $Y \in W^{1,1}(0,\infty)$, $\alpha \in L^1_{loc}(0,\infty)$, and T > 0 be fixed. Suppose that

$$\frac{\mathrm{d}}{\mathrm{d}t}Y(t) + \alpha(t)Y(t) \le 0 \qquad a.e. \ in \ (0, \infty),$$

where

$$\int_t^{t+T} \alpha(s) \, ds \geq \beta > 0, \qquad \quad \int_t^{t+T} \, \alpha^-(s) \, ds \leq M < \infty, \qquad \quad \forall \, t \geq t_0,$$

with $\alpha^{-}(s) = \max\{-\alpha(s), 0\}$, and $t_0 \geq 0$. Then, we have

$$Y(t) \le Y(t_0) e^{(\beta+M)} e^{-\left(\frac{t-t_0}{T}\right)\beta}, \quad \forall t \ge t_0.$$

Proof. An application of Lemma 2.3 gives

$$Y(t) \le Y(t_0) e^{-\int_{t_0}^t \alpha(s) ds}, \qquad \forall t \ge t_0.$$

For any $t \geq t_0$, there exists $K \in \mathbb{N}$ such that $t_0 + KT \leq t \leq t_0 + (K+1)T$. Then, we have

$$e^{-\int_{t_0}^t \alpha(s) \, ds} = \exp\left(-\int_{t_0}^{t_0+T} \alpha(s) \, ds\right) \cdots \exp\left(-\int_{t_0+(K-1)T}^{t_0+KT} \alpha(s) \, ds\right) \exp\left(-\int_{t_0+KT}^t \alpha(s) \, ds\right)$$

$$\leq e^{-K\beta} \exp\left(\int_{t_0+KT}^{t_0+(K+1)T} \alpha^-(s) \, ds\right) \leq e^{-\left(\frac{t-t_0}{T}-1\right)\beta} e^{M} \leq e^{-\left(\frac{t-t_0}{T}\right)\beta} e^{\beta+M}.$$

The proof is complete. \square

3. The Ladyzhenskaya model

The phenomenon that we consider in this section is the motion of an incompressible viscous fluid in a bounded Lipschitz domain $\Omega \subset \mathbb{R}^d$, $d \in \{2,3\}$ with no-slip boundary conditions. Let **u** denote the velocity field, P the pressure, and **f** the body force per unit mass. In [35], Ladyzhenskaya proposed the following mathematical model

$$\partial_{t}\mathbf{u} + (\mathbf{u} \cdot \nabla)\mathbf{u} - \nabla \cdot \mathbf{T}(\mathbf{D}\mathbf{u}) + \nabla P = \mathbf{f},$$

$$\nabla \cdot \mathbf{u} = 0,$$

$$\mathbf{u}|_{\partial O} = 0.$$
(3.1)

where \mathbf{T} denotes the Cauchy stress of an incompressible and homogeneous fluid whose constitutive relation is given by

$$\mathbf{T}(\mathrm{D}\mathbf{u}) = 2\left(\nu_0 + \nu_1 |\mathrm{D}\mathbf{u}|_F^{p-2}\right) \mathrm{D}\mathbf{u}, \quad p \ge 2, \tag{3.2}$$

with initial condition $\mathbf{u}(\cdot,0) = \mathbf{u}_0(\cdot)$. Here, $D\mathbf{u} = \frac{1}{2} \left[(\nabla \mathbf{u}) + (\nabla \mathbf{u})^T \right]$, and ν_0 and ν_1 are positive parameters. It is worth mentioning that ν_0 scales as $\frac{(\text{length})^2}{(\text{time})}$, and ν_1 has dimension $(\text{time})^{p-3} \times (\text{length})^2$. In the literature, some works have been devoted to the case with $\mathbf{T} = \mathbf{T}(\nabla \mathbf{u})$, namely $D\mathbf{u}$ is replaced by the full velocity gradient $\nabla \mathbf{u}$ in (3.2). However, in such a case the model does not comply with the principle of frame indifference (see, e.g., [38,39]).

Before stating the well-posedness result, we report the following property of the constitutive relation (3.2), which will be of key usefulness in the sequel. In particular, we will exploit the factor ν_0 . We refer the reader to [38,49] for the proof.

Proposition 3.1. Let **T** be as given in (3.2). For all $A, B \in \mathbb{R}^{3\times 3}_{sum}$, we have

$$(\mathbf{T}(A) - \mathbf{T}(B)) : (A - B) \ge 2 \nu_0 |A - B|_F^2.$$
 (3.3)

Taking advantage of the enhanced regularity due to (3.2), Ladyzhenskaya showed in [35,49] that the weak solutions to (3.1) are global in time and unique for any Reynolds number and any exponent $p \geq \frac{5}{2}$. For an overview, we refer the reader to [36] and [37, Theorem 7.2], and to [37,65,66] for the existence of compact finite dimensional global attractor. Later on, many contributions have been devoted to the analysis of the case $1 \leq p < \frac{5}{2}$. Without any claim to give an exhaustive survey, we mention the existence of global measure-valued solutions for $\frac{6}{5} , global weak solutions for <math>p > \frac{9}{5}$ and global strong solutions for $p \geq \frac{9}{4}$ obtained in [38,52,53]. For the periodic case, enhanced results in terms of p have been achieved as reported in [39,40]. In particular, the existence, but not uniqueness, of global weak solutions fulfilling the energy equality holds for $p \geq \frac{11}{5}$. Moreover, under additional assumptions on the initial datum and the forcing term, global in time and unique strong solutions also exist. The asymptotic behavior in the same range of p has been studied in [51,67].

We now state the well-posedness result for the model (3.1) proved in [35] (see also [38]).

Theorem 3.2 (Existence and Uniqueness of Weak Solutions). Assume that $p \geq \frac{5}{2}$, $\mathbf{f} \in L^2(0,T;L^2(\Omega))$ and $\mathbf{u}_0 \in L^2_{\sigma}(\Omega)$. Problem (3.1) has a unique weak solution on $(0,\infty)$ satisfying for all T > 0

$$\mathbf{u} \in \mathcal{C}([0,T]; L^2_{\sigma}(\Omega)) \cap L^p(0,T; W^{1,p}_{\sigma}(\Omega)), \quad \partial_t \mathbf{u} \in L^{p'}(0,T; (W^{1,p}_{\sigma}(\Omega))'),$$

where p' is the conjugate exponent of p, and

$$\langle \partial_t \mathbf{u}, \mathbf{w} \rangle + ((\mathbf{u} \cdot \nabla)\mathbf{u}, \mathbf{w}) + (\mathbf{T}(\mathrm{D}\mathbf{u}), \nabla \mathbf{w}) = (\mathbf{f}, \mathbf{w}), \quad \forall \mathbf{w} \in W^{1,p}_{\sigma}(\Omega),$$

for almost all $t \in [0,T]$. Moreover, the energy equality holds

$$\frac{1}{2}\|\mathbf{u}(t)\|^{2} + \int_{0}^{t} \left(2\nu_{0}\|\mathbf{D}\mathbf{u}(\tau)\|^{2} + 2\nu_{1}\|\mathbf{D}\mathbf{u}(\tau)\|_{L^{p}}^{p}\right) d\tau = \frac{1}{2}\|\mathbf{u}_{0}\|^{2} + \int_{0}^{t} (\mathbf{f}(\tau), \mathbf{u}(\tau)) d\tau, \qquad \forall t \ge 0.$$
 (3.4)

Let λ_1 be the smallest eigenvalue of the Stokes operator. Assume that \mathbf{f} is time independent. We denote by G the Grashof number in three-dimensions defined as

$$G = \frac{\|\mathbf{f}\|}{\nu_0^2 \lambda_1^{3/4}}. (3.5)$$

We now give bounds on the solution \mathbf{u} of (3.1) that will be used in our analysis.

Proposition 3.3. Fix T > 0, and let $\mathbf{f} \in L^2(\Omega)$. Suppose that \mathbf{u} is a weak solution of (3.1), then we have

$$\|\mathbf{u}(t)\|^{2} \leq \|\mathbf{u}_{0}\|^{2} e^{-\lambda_{1}\nu_{0}t} + \frac{\|\mathbf{f}\|^{2}}{\lambda_{1}^{2}\nu_{0}^{2}} \left(1 - e^{-\lambda_{1}\nu_{0}t}\right), \qquad \forall t \geq 0.$$
(3.6)

As a consequence, there exists a time $t_0 > 0$ such that for all $t \ge t_0$ we have

$$\|\mathbf{u}(t)\|^2 \le 2\frac{\nu_0^2 G^2}{\lambda_1^{\frac{1}{2}}} \tag{3.7}$$

and

$$\int_{t}^{t+T} \left(\nu_{0} \| \mathbf{D} \mathbf{u}(\tau) \|^{2} + \nu_{1} \| \mathbf{D} \mathbf{u}(\tau) \|_{L^{p}}^{p} \right) d\tau \leq 2 \left(1 + \nu_{0} \lambda_{1} T \right) \frac{\nu_{0}^{2} G^{2}}{\lambda_{1}^{\frac{1}{2}}}.$$
(3.8)

The proof of Proposition 3.3 is standard and thus omitted here. For the readers' convenience, we observe that (3.6) follows from Lemma 2.3 after dropping the ν_1 term in (3.4). In addition, (3.8) is a consequence of (3.4) and (3.7).

4. The case $p \geq \frac{5}{2}$ with no-slip boundary conditions

In this section, we first analyze the nudging algorithm for the Ladyzhenskaya model with no-slip boundary conditions for $p \geq \frac{5}{2}$. After proving the global well-posedness of the week solution in Theorem 4.1, we proceed to the task of finding conditions on h and μ under which the approximate solution obtained by this algorithm converges to the reference solution over time, summarized in Theorem 4.2.

Let $I_h(\mathbf{u}(t))$ represent the observational measurements at a spatial resolution of size h for t > 0 satisfying (2.1). The approximating solution \mathbf{v} with initial condition $\mathbf{v}(\cdot,0) = \mathbf{v}_0(\cdot)$, chosen arbitrarily, shall be given by

$$\partial_{t}\mathbf{v} + (\mathbf{v} \cdot \nabla)\mathbf{v} - \nabla \cdot \mathbf{T}(\mathbf{D}\mathbf{v}) + \nabla Q = \mathbf{f} - \mu I_{h}(\mathbf{v} - \mathbf{u}),$$

$$\nabla \cdot \mathbf{v} = 0,$$

$$\mathbf{v}|_{\partial Q} = 0.$$
(4.1)

The first result of this manuscript concerns the global well-posedness of weak solutions for the Data Assimilation algorithm.

Theorem 4.1. Assume that $p \geq \frac{5}{2}$, $\mathbf{f} \in L^2(0,T;L^2(\Omega))$ and $\mathbf{v}_0 \in L^2_{\sigma}(\Omega)$. Let \mathbf{u} be the solution to problem (3.1) from Theorem 3.2. The continuous data assimilation Eqs. (4.1) has a unique global weak solution that satisfies for all T > 0

$$\mathbf{v} \in \mathcal{C}([0,T]; L^2_{\sigma}(\Omega)) \cap L^p(0,T; W^{1,p}_{\sigma}(\Omega)) \cap W^{1,p'}(0,T; (W^{1,p}_{\sigma}(\Omega))')$$

where $p' = \frac{p}{p-1}$, and

$$\langle \partial_t \mathbf{v}, \mathbf{w} \rangle + ((\mathbf{v} \cdot \nabla) \mathbf{v}, \mathbf{w}) + (\mathbf{T}(\mathbf{D}\mathbf{v}), \nabla \mathbf{w}) = (\mathbf{f}, \mathbf{w}) - \mu(I_h(\mathbf{v} - \mathbf{u}), \mathbf{w}), \quad \forall \mathbf{w} \in W_{\sigma}^{1,p}(\Omega),$$
 (4.2)

for almost all $t \in [0, T]$.

Proof. The strategy is to reformulate (4.1) as a fixed-point problem. For any fixed T > 0, define

$$\mathcal{F}: L^{2}(0,T;L^{2}_{\sigma}(\varOmega)) \to \mathcal{C}([0,T];L^{2}_{\sigma}(\varOmega)) \cap L^{p}(0,T;W^{1,p}_{\sigma}(\varOmega)) \cap W^{1,p'}(0,T;(W^{1,p}_{\sigma}(\varOmega))') \tag{4.3}$$

by

$$\mathcal{F}(\mathbf{w}) = \mathbf{v},$$

where \mathbf{v} is a weak solution to the problem

$$\partial_{t}\mathbf{v} + (\mathbf{v} \cdot \nabla)\mathbf{v} - \nabla \cdot \mathbf{T}(\mathbf{D}\mathbf{v}) + \nabla q = \mathbf{f}_{\mu} - \mu I_{h}\mathbf{w},$$

$$\nabla \cdot \mathbf{v} = 0,$$

$$\mathbf{v}|_{\partial\Omega} = 0,$$

$$\mathbf{v}(\cdot, 0) = \mathbf{v}_{0}(\cdot),$$

$$(4.4)$$

for a given $\mathbf{w} \in L^2(0,T;L^2_{\sigma}(\Omega))$, with $\mathbf{f}_{\mu} = \mathbf{f} + \mu I_h \mathbf{u}$. It is easy to verify that $\mathbf{f}_{\mu} \in L^2(0,T;L^2(\Omega))$ since I_h is a continuous and bounded linear operator. The above map \mathcal{F} is well-defined since the existence and uniqueness of a weak solution \mathbf{v} for any given initial condition $\mathbf{v}_0 \in L^2_{\sigma}(\Omega)$ follows directly from Theorem 3.2. Now, define

$$\mathcal{A} = \left\{ \mathbf{w} \in L^2(0, T; L^2_{\sigma}(\Omega)) : \int_0^t \|\mathbf{w}(\tau)\|^2 d\tau \le c_1 e^{c_2 t}, \quad \forall t \in [0, T] \right\}, \tag{4.5}$$

with

$$c_1 = 2T \left(\|\mathbf{v}_0\|_{L^2_{\sigma}} + \int_0^T \|\mathbf{f}_{\mu}(\tau)\| d\tau \right)^2, \quad c_2 = 2 \,\mu^2 \,c_I^2 \,T.$$

To apply the Schauder fixed-point theorem (see Theorem 2.1) to the above problem, we will verify the theorem's assumptions in the next five steps.

Step I. We claim that $\mathcal{F}: \mathcal{A} \to \mathcal{A}$, i.e. $\mathcal{F}(\mathcal{A}) \subset \mathcal{A}$.

From the energy equality (2.1) and (3.4), we have

$$\|\mathbf{v}(\tau)\| \le \|\mathbf{v}_0\| + \int_0^{\tau} \|\mathbf{f}_{\mu}(s)\| \, ds + \mu \int_0^{\tau} \|I_h \mathbf{w}(s)\| \, ds$$

$$\le \|\mathbf{v}_0\| + \int_0^{\tau} \|\mathbf{f}_{\mu}(s)\| \, ds + \mu \, c_I \int_0^{\tau} \|\mathbf{w}(s)\| \, ds$$
(4.6)

for all $\tau \in [0, T]$. By using Young's inequality and the Hölder's inequality, we obtain

$$\|\mathbf{v}(\tau)\|^{2} \leq 2\left(\|\mathbf{v}_{0}\| + \int_{0}^{\tau} \|\mathbf{f}_{\mu}(s)\| ds\right)^{2} + 2\mu^{2} c_{I}^{2} \left(\int_{0}^{\tau} \|\mathbf{w}(s)\| ds\right)^{2}$$

$$\leq 2\left(\|\mathbf{v}_{0}\| + \int_{0}^{\tau} \|\mathbf{f}_{\mu}(s)\| ds\right)^{2} + 2\mu^{2} c_{I}^{2} T \left(\int_{0}^{\tau} \|\mathbf{w}(s)\|^{2} ds\right).$$
(4.7)

Since $\mathbf{w} \in \mathcal{A}$, we infer that

$$\int_{0}^{t} \|\mathbf{v}(\tau)\|^{2} d\tau \leq 2 \int_{0}^{t} \left(\|\mathbf{v}_{0}\| + \int_{0}^{T} \|\mathbf{f}_{\mu}(s)\| ds \right)^{2} d\tau + 2 \mu^{2} c_{I}^{2} T \int_{0}^{t} \int_{0}^{\tau} \|\mathbf{w}(s)\|^{2} ds d\tau
\leq c_{1} + c_{2} \int_{0}^{t} \left(c_{1} e^{c_{2}\tau} \right) d\tau = c_{1} e^{c_{2}t},$$
(4.8)

which, in turn, entails $\mathbf{v} = \mathcal{F}(\mathbf{w}) \in \mathcal{A}$.

Step II. \mathcal{A} is a closed set in $L^2(0,T;L^2_{\sigma}(\Omega))$.

Assume that $\{\mathbf{w}_n\}_{n=0}^{\infty} \subset \mathcal{A}$ is such that $\mathbf{w}_n \to \mathbf{w}$ in $L^2(0,T;L^2_{\sigma}(\Omega))$. It follows that \mathcal{A} is closed from the following argument

$$\int_0^t \|\mathbf{w}(\tau)\|^2 d\tau = \lim_{n \to \infty} \int_0^t \|\mathbf{w}_n(\tau)\|^2 d\tau \le \lim_{n \to \infty} c_1 e^{c_2 t} = c_1 e^{c_2 t}, \quad \forall t \in [0, T].$$

Step III. \mathcal{A} is convex set in $L^2(0,T;L^2_{\sigma}(\Omega))$.

Let $\mathbf{w}_1, \mathbf{w}_2 \in \mathcal{A}$, then $\lambda \mathbf{w}_1 + (1 - \lambda) \mathbf{w}_2 \in L^2 \left(0, T; L^2_{\sigma}(\Omega)\right)$ for any $\lambda \in [0, 1]$. We compute

$$\begin{split} & \int_{0}^{t} \|\lambda \mathbf{w}_{1}(\tau) + (1 - \lambda)\mathbf{w}_{2}(\tau)\|^{2} d\tau \\ & = \lambda^{2} \int_{0}^{t} \|\mathbf{w}_{1}(\tau)\|^{2} d\tau + 2\lambda(1 - \lambda) \int_{0}^{t} (\mathbf{w}_{1}(\tau), \mathbf{w}_{2}(\tau)) d\tau + (1 - \lambda)^{2} \int_{0}^{t} \|\mathbf{w}_{2}(\tau)\|^{2} d\tau \\ & \leq \lambda^{2} \int_{0}^{t} \|\mathbf{w}_{1}(\tau)\|^{2} d\tau + 2\lambda(1 - \lambda) \left(\int_{0}^{t} \|\mathbf{w}_{1}(\tau)\|^{2} d\tau \right)^{\frac{1}{2}} \left(\int_{0}^{t} \|\mathbf{w}_{2}(\tau)\|^{2} d\tau \right)^{\frac{1}{2}} \\ & + (1 - \lambda)^{2} \int_{0}^{t} \|\mathbf{w}_{2}(\tau)\|^{2} d\tau \\ & \leq \left(\lambda^{2} + (1 - \lambda)^{2} + 2\lambda(1 - \lambda)\right) c_{1}e^{c_{2}t} \\ & = c_{1}e^{c_{2}t}, \end{split}$$

which means $\lambda \mathbf{w}_1 + (1 - \lambda) \mathbf{w}_2 \in \mathcal{A}$, proving the convexity.

Step IV. $\mathcal{F}: \mathcal{A} \to \mathcal{A}$ is continuous.

Consider $\{\mathbf{w}_n\}_{n=1}^{\infty} \subset \mathcal{A}$ such that $\mathbf{w}_n \to \mathbf{w}$ in $L^2(0,T;L^2_{\sigma}(\Omega))$. We are required to show that $\mathbf{v}_n = \mathcal{F}(\mathbf{w}_n) \to \mathcal{F}(\mathbf{w}) = \mathbf{v}$ in $L^2(0,T;L^2_{\sigma}(\Omega))$. First, define the difference $\psi_n = \mathbf{v}_n - \mathbf{v}$, which solves

$$\langle \partial_t \psi_n, \varphi \rangle + (\mathbf{v}_n \cdot \nabla \mathbf{v}_n, \varphi) - (\mathbf{v} \cdot \nabla \mathbf{v}, \varphi) + (\mathbf{T}(\mathbf{D}\mathbf{v}_n) - \mathbf{T}(\mathbf{D}\mathbf{v}), \nabla \varphi) = -\mu (I_h(\mathbf{w}_n - \mathbf{w}), \varphi)$$

for all $\varphi \in W_{\sigma}^{1,3}(\Omega)$, for almost all $t \in [0,T]$. Thanks to [38, Lemma 2.45], the incompressibility condition and the regularity (4.3), choosing $\varphi = \psi_n$ in the above equation, we obtain

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{dt}} \|\psi_n\|^2 + (\psi_n \cdot \nabla \mathbf{v}, \psi_n) + (\mathbf{T}(\mathrm{D}\mathbf{v}_n) - \mathbf{T}(\mathrm{D}\mathbf{v}), \mathrm{D}\mathbf{v}_n - \mathrm{D}\mathbf{v}) = -\mu \left(I_h(\mathbf{w}_n - \mathbf{w}), \psi_n\right).$$

By exploiting (3.3), the Korn inequality, the (Hölder inequality) with $p' = \frac{p}{p-1}$ and the (Lebesgue interpolation inequality) in L^p -spaces with $\theta = 1 - \frac{3}{2p}$, we find

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{dt}} \|\psi_{n}\|^{2} + \nu_{0} \|\nabla\psi_{n}\|^{2} \leq \left| (\psi_{n} \cdot \nabla \mathbf{v}, \psi_{n}) \right| + \left| \mu \left(I_{h}(\mathbf{w}_{n} - \mathbf{w}), \psi_{n} \right) \right| \\
\leq \|\psi_{n}\|_{L^{2p'}}^{2} \|\nabla \mathbf{v}\|_{L^{p}} + \mu \|I_{h}(\mathbf{w}_{n} - \mathbf{w})\| \|\psi_{n}\| \\
\leq \|\psi_{n}\|^{2 - \frac{3}{p}} \|\psi_{n}\|_{L^{6}}^{\frac{3}{p}} \|\nabla \mathbf{v}\|_{L^{p}} + \mu c_{I} \|\mathbf{w}_{n} - \mathbf{w}\| \|\psi_{n}\| \\
\leq c_{S} \|\psi_{n}\|^{2 - \frac{3}{p}} \|\nabla\psi_{n}\|^{\frac{3}{p}} \|\nabla \mathbf{v}\|_{L^{p}} + \mu c_{I} \|\mathbf{w}_{n} - \mathbf{w}\| \|\psi_{n}\| \\
\leq c_{S} \|\psi_{n}\|^{2 - \frac{3}{p}} \|\nabla\psi_{n}\|^{\frac{3}{p}} \|\nabla \mathbf{v}\|_{L^{p}} + \mu c_{I} \|\mathbf{w}_{n} - \mathbf{w}\| \|\psi_{n}\| \\
\leq \frac{\nu_{0}}{2} \|\nabla\psi_{n}\|^{2} + \left(\tilde{c} \nu_{0}^{-\frac{3}{2p-3}} c_{S}^{\frac{2p}{2p-3}} \|\nabla \mathbf{v}\|_{L^{p}}^{\frac{2p}{2p-3}} + \frac{1}{4}\right) \|\psi_{n}\|^{2} + \mu^{2} c_{I}^{2} \|\mathbf{w}_{n} - \mathbf{w}\|^{2},$$

where \tilde{c} only depends on p. In the above estimate, the constant c_S denotes the (Sobolev embedding) $H_0^1(\Omega) \hookrightarrow L^6(\Omega)$. Therefore, we obtain

$$\frac{\mathrm{d}}{\mathrm{d}t} \|\psi_n\|^2 \le \left(\frac{1}{4} + \tilde{c}\,\nu_0^{-\frac{3}{2p-3}}\,c_S^{\frac{2p}{2p-3}} \|\nabla\mathbf{v}\|_{L^p}^{\frac{2p}{2p-3}}\right) \|\psi_n\|^2 + \mu^2\,c_I^2 \|\mathbf{w}_n - \mathbf{w}\|^2.$$

Applying the Gronwall lemma (see Lemma 2.3) to the above inequality, we get

$$\|\psi_n(t)\|^2 \le \|\psi_n(0)\|^2 e^{\int_0^t \lambda(\tau) d\tau} + \mu^2 c_I^2 \int_0^t \|\mathbf{w}_n(s) - \mathbf{w}(s)\|^2 e^{\int_s^t \lambda(\tau) d\tau} ds,$$

for all $t \in [0, T]$, where

$$\lambda(\tau) = \frac{1}{4} + \tilde{c} \,\nu_0^{-\frac{3}{2p-3}} \, c_S^{\frac{2p}{2p-3}} \, \|\nabla \mathbf{v}(\tau)\|_{L^p}^{\frac{2p}{2p-3}}.$$

Note that having $p \geq \frac{5}{2}$ yields $\frac{2p}{2p-3} \leq p$, thereby the regularity $\mathbf{v} \in L^p(0,T;W^{1,p}_\sigma(\Omega))$ entails that $\lambda(\tau) \in L^1[0,T]$. In light of $\psi_n(0) = 0$, we are led to

$$\|\psi_n\|_{L^{\infty}(0,T;L^2_{\sigma}(\Omega))} \le \mu \, c_I \, e^{\frac{1}{2}\|\lambda\|_{L^1(0,T)}} \, \|\mathbf{w}_n - \mathbf{w}\|_{L^2(0,T;L^2_{\sigma}(\Omega))}.$$

Since the right-hand side converges to 0 as $n \to \infty$, this implies the continuity of \mathcal{F} .

Step V. We construct a compact subset \mathcal{K} of \mathcal{A} such that $\mathcal{F}(\mathcal{A}) \subset \mathcal{K}$. From the energy equality (3.4) written for the solution to (4.4), and after using the (Hölder inequality), the Korn inequality and (2.1), we have

$$\|\mathbf{v}(t)\|^{2} + \int_{0}^{t} \left(2\nu_{0}\|\nabla\mathbf{v}(\tau)\|^{2} + \frac{4\nu_{1}}{c_{K}^{p}}\|\nabla\mathbf{v}(\tau)\|_{L^{p}}^{p}\right) d\tau \leq \|\mathbf{v}_{0}\|^{2} + 2\int_{0}^{t} \left(\|\mathbf{f}_{\mu}(\tau)\| + \mu\|I_{h}\mathbf{w}(\tau)\|\right)\|\mathbf{v}(\tau)\| d\tau$$

$$\leq \|\mathbf{v}_{0}\|^{2} + \frac{2}{\lambda_{1}} \int_{0}^{t} \left(\|\mathbf{f}_{\mu}(\tau)\| + \mu\|I_{h}\mathbf{w}(\tau)\|\right)\|\nabla\mathbf{v}(\tau)\| d\tau$$

$$\leq \|\mathbf{v}_{0}\|^{2} + \nu_{0} \int_{0}^{t} \|\nabla\mathbf{v}(\tau)\|^{2} d\tau + \frac{1}{2\nu_{0}\lambda_{1}^{2}} \int_{0}^{t} \left(\|\mathbf{f}_{\mu}(\tau)\| + \mu\|I_{h}\mathbf{w}(\tau)\|\right)^{2} d\tau$$

$$\leq \|\mathbf{v}_{0}\|^{2} + \nu_{0} \int_{0}^{t} \|\nabla\mathbf{v}(\tau)\|^{2} d\tau + \frac{1}{\nu_{0}\lambda_{1}^{2}} \int_{0}^{t} \left(\|\mathbf{f}_{\mu}(\tau)\|^{2} + \mu^{2}c_{1}^{2}\|\mathbf{w}(\tau)\|^{2}\right) d\tau,$$

for all $t \in [0, T]$. Thus, we arrive at

$$\|\mathbf{v}(t)\|^2 + \int_0^t \left(\nu_0 \|\nabla \mathbf{v}(\tau)\|^2 + \frac{4\nu_1}{c_{\mathrm{K}}^p} \|\nabla \mathbf{v}(\tau)\|_{L^p}^p\right) \, d\tau \leq \|\mathbf{v}_0\|^2 + \frac{1}{\nu_0 \lambda_1^2} \|\mathbf{f}_\mu\|_{L^2(0,T;L^2_\sigma(\varOmega))}^2 \, + \, \frac{\mu^2 c_1^2}{\nu_0 \lambda_1^2} \, c_1 e^{c_2 T} := \tilde{c}_0,$$

With \tilde{c}_0 defined as above, we deduce that

$$\|\mathbf{v}\|_{L^{\infty}(0,T;L^{2}_{\sigma}(\Omega))} \leq \sqrt{\tilde{c}_{0}} := \tilde{c}_{1}, \quad \|\mathbf{v}\|_{L^{p}(0,T;W^{1,p}_{\sigma}(\Omega))} \leq \left(\frac{\tilde{c}_{0}c_{K}^{p}}{2\nu_{1}}\right)^{\frac{1}{p}} := \tilde{c}_{2}. \tag{4.9}$$

Then, we infer that

$$\mathcal{F}(\mathcal{A}) \subset \mathcal{B} = \left\{ \mathbf{v} \in \mathcal{A} : \|\mathbf{v}\|_{L^{\infty}\left(0,T;L^{2}_{\sigma}(\Omega)\right)} \leq \tilde{c}_{1} \quad \text{and} \quad \|\mathbf{v}\|_{L^{p}\left(0,T;W^{1,p}_{\sigma}(\Omega)\right)} \leq \tilde{c}_{2} \right\}.$$

Next we investigate the time derivative $\partial_t \mathbf{v}$. We recall the weak formulation of (4.4)

$$\langle \partial_t \mathbf{v}, \varphi \rangle + \int_{\Omega} \mathbf{v} \cdot \nabla \mathbf{v} \cdot \varphi \, d\mathbf{x} + \int_{\Omega} 2\nu_0 \, \mathrm{D} \mathbf{v} : \mathrm{D} \varphi + 2\nu_1 \, |\mathrm{D} \mathbf{v}|_F^{p-2} \, \mathrm{D} \mathbf{v} : \mathrm{D} \varphi \, d\mathbf{x}$$
$$= \int_{\Omega} \mathbf{f}_{\mu} \cdot \varphi \, d\mathbf{x} - \mu \int_{\Omega} I_h \mathbf{w} \cdot \varphi \, d\mathbf{x}$$

for all $\varphi \in W^{1,p}_{\sigma}(\Omega)$, for almost all $t \in [0,T]$. Due to the incompressibility condition, the nonlinear term can be written as $\mathbf{v} \cdot \nabla \mathbf{v} = \nabla \cdot (\mathbf{v} \otimes \mathbf{v})$. Then, we have

$$|\langle \partial_t \mathbf{v}, \varphi \rangle| \leq \left| \int_{\Omega} \mathbf{v} \otimes \mathbf{v} : \nabla \varphi \, d\mathbf{x} \right| + \left| \int_{\Omega} 2\nu_0 \, \mathrm{D} \mathbf{v} : \nabla \varphi + 2\nu_1 \, |\mathrm{D} \mathbf{v}|_F^{p-2} \, \mathrm{D} \mathbf{v} : \nabla \varphi \, d\mathbf{x} \right| + \left| \int_{\Omega} \mathbf{f}_{\mu} \cdot \varphi \, d\mathbf{x} \right| + \mu \left| \int_{\Omega} I_h \mathbf{w} \cdot \varphi \, d\mathbf{x} \right|.$$

Let $p' = \frac{p}{p-1}$, and note that p' < p for $p \ge \frac{5}{2}$. Using the (Hölder inequality) along with (2.1) yields

$$|\langle \partial_t \mathbf{v}, \varphi \rangle| \leq \|\mathbf{v}\|_{L^{2p'}}^2 \|\nabla \varphi\|_{L^p} + 2\nu_0 \|\mathbf{D}\mathbf{v}\|_{L^{p'}} \|\nabla \varphi\|_{L^p} + 2\nu_1 \|\mathbf{D}\mathbf{v}\|_{L^p}^{p-1} \|\nabla \varphi\|_{L^p} + \|\mathbf{f}_{\mu}\| \|\varphi\| + \mu c_I \|\mathbf{w}\| \|\varphi\|.$$

By taking supremum of the above inequality over all $\varphi \in W^{1,p}_{\sigma}(\Omega)$ such that $\|\varphi\|_{W^{1,p}_{\sigma}(\Omega)} = 1$, and using the (Lebesgue interpolation inequality), we obtain

$$\begin{split} \|\partial_{t}\mathbf{v}\|_{\left(W_{\sigma}^{1,p}(\Omega)\right)'} &\leq \|\mathbf{v}\|_{L^{2p'}}^{2} + 2\nu_{0}\|\mathbf{D}\mathbf{v}\|_{L^{p'}} + 2\nu_{1}\|\mathbf{D}\mathbf{v}\|_{L^{p}}^{p-1} + C\|\mathbf{f}_{\mu}\| + \mu c_{I} C\|\mathbf{w}\| \\ &\leq C\|\mathbf{v}\|^{\frac{2p-3}{p}} \|\nabla\mathbf{v}\|^{\frac{3}{p}} + \nu_{0} C\|\nabla\mathbf{v}\|_{L^{2}} + \nu_{1} C\|\nabla\mathbf{v}\|_{L^{p}}^{p-1} + C\|\mathbf{f}_{\mu}\| + \mu c_{I} C\|\mathbf{w}\| \\ &\leq C\|\mathbf{v}\|^{\frac{2p-3}{p}} \|\nabla\mathbf{v}\|_{L^{p}}^{\frac{3}{p}} + \nu_{0} C\|\nabla\mathbf{v}\|_{L^{p}} + \nu_{1} C\|\nabla\mathbf{v}\|_{L^{p}}^{p-1} + C\|\mathbf{f}_{\mu}\| + \mu c_{I} C\|\mathbf{w}\|, \end{split}$$

where C only depends on p and Ω . Hence,

$$\begin{split} &\|\partial_{t}\mathbf{v}\|_{L^{p'}\left(0,T;\left(W_{\sigma}^{1,p}(\Omega)\right)'\right)}^{p'} = \int_{0}^{T} \|\partial_{t}\mathbf{v}(\tau)\|_{\left(W_{\sigma}^{1,p}(\Omega)\right)'}^{p'} d\tau \\ &\leq C \int_{0}^{T} \|\mathbf{v}(\tau)\|_{L^{p}}^{\frac{2p-3}{p-1}} \|\nabla\mathbf{v}(\tau)\|_{L^{p}}^{\frac{3}{p-1}} d\tau + \nu_{0}^{p'} C \int_{0}^{T} \|\nabla\mathbf{v}(\tau)\|_{L^{p}}^{\frac{p}{p-1}} d\tau + \nu_{1}^{p'} C \int_{0}^{T} \|\nabla\mathbf{v}(\tau)\|_{L^{p}}^{p} d\tau \\ &+ C \int_{0}^{T} \|\mathbf{f}_{\mu}(\tau)\|_{p'}^{p'} d\tau + \mu^{p'} c_{I}^{p'} C \int_{0}^{T} \|\mathbf{w}(\tau)\|_{p'}^{p'} dt \\ &\leq C \|\mathbf{v}\|_{L^{\infty}\left(0,T;L_{\sigma}^{2}(\Omega)\right)}^{\frac{2p-3}{p-1}} T^{\frac{1}{\alpha}} \|\mathbf{v}\|_{L^{p}\left(0,T;W_{\sigma}^{1,p}(\Omega)\right)}^{\frac{3}{p-1}} + \nu_{0}^{p'} C T^{\frac{1}{\beta}} \|\mathbf{v}\|_{L^{p}\left(0,T;W_{\sigma}^{1,p}(\Omega)\right)}^{p'} \\ &+ \nu_{1}^{p'} C \|\mathbf{v}\|_{L^{p}\left(0,T;W_{\sigma}^{1,p}(\Omega)\right)}^{p} + \mu^{p'} c_{I}^{p'} C T^{\frac{1}{\gamma}} \|\mathbf{w}\|_{L^{2}\left(0,T;L_{\sigma}^{2}(\Omega)\right)}^{p'} \\ &+ C T^{\frac{1}{\gamma}} \|\mathbf{f}_{\mu}\|_{L^{2}\left(0,T;L_{\sigma}^{2}(\Omega)\right)}^{p'} + \mu^{p'} c_{I}^{p'} C T^{\frac{1}{\gamma}} \|\mathbf{w}\|_{L^{2}\left(0,T;L_{\sigma}^{2}(\Omega)\right)}^{p'} \\ &\coloneqq \tilde{c}_{3}^{p'}, \end{split}$$

where α, β and γ are the conjugate exponents to (p-1)p/3, p-1 and 2/p', respectively, and the constant C depends only on p and Ω . Given $\tilde{c}_3^{p'}$ as above, we have

$$\|\partial_t \mathbf{v}\|_{L^{p'}\left(0,T;\left(W_{\sigma}^{1,p}(\Omega)\right)'\right)} \le \tilde{c}_3. \tag{4.10}$$

Finally, with $\tilde{c}_1, \tilde{c}_2, \tilde{c}_3$ given in (4.9) and (4.10), respectively, we infer that

$$\mathcal{F}(\mathcal{A}) \subset \mathcal{K},$$
 (4.11)

where

$$\mathcal{K} = \left\{ \mathbf{v} \in \mathcal{A} : \|\mathbf{v}\|_{L^{\infty}\left(0,T;L^{2}_{\sigma}(\Omega)\right)} \leq \tilde{c}_{1}, \ \|\mathbf{v}\|_{L^{p}\left(0,T;W^{1,p}_{\sigma}(\Omega)\right)} \leq \tilde{c}_{2} \text{ and } \|\partial_{t}\mathbf{v}\|_{L^{p'}\left(0,T;\left(W^{1,p}_{\sigma}(\Omega)\right)'\right)} \leq \tilde{c}_{3} \right\}.$$

We are left to show that \mathcal{K} is a compact subset of \mathcal{A} . Since $W_{\sigma}^{1,p}(\Omega) \subset L_{\sigma}^{2}(\Omega) \subset \left(W_{\sigma}^{1,p}(\Omega)\right)'$, thanks to Theorem 2.2, we deduce that \mathcal{K} is compactly embedded in $L^{p}\left(0,T,L_{\sigma}^{2}(\Omega)\right)$, and, in turn, in $L^{2}\left(0,T,L_{\sigma}^{2}(\Omega)\right)$ since $p \geq \frac{5}{2}$. Therefore, to summarize it is proved that

$$\mathcal{F}(\mathcal{A}) \subset \mathcal{K} \stackrel{c}{\hookrightarrow} \mathcal{A}$$

where \mathcal{K} is a compact subset of \mathcal{A} with respect to the norm $L^2\left(0,T,L^2_{\sigma}(\Omega)\right)$. As a consequence of Theorem 2.1, $\mathcal{F}:\mathcal{A}\to\mathcal{A}$ has a fixed point in \mathcal{K} , which implies the existence result in Theorem 4.1. Lastly, the uniqueness of the weak solution to problem (4.1) is obtained from the same argument of Step IV by replacing \mathbf{v}_n and \mathbf{v} with two solutions \mathbf{v}_1 and \mathbf{v}_2 , respectively, originating from the same initial datum \mathbf{v}_0 . \square

Next, we prove the convergence result.

Theorem 4.2. For $p \geq \frac{5}{2}$, let $\mathbf{f} \in L^2(\Omega)$ and let \mathbf{u} be a weak solution of (3.1) with no-slip Dirichlet boundary conditions departing from $\mathbf{u}_0 \in L^2_{\sigma}(\Omega)$. Let \mathbf{v} be the solution to the data assimilation algorithm given by (4.1). Then, for μ large enough such that

$$\mu \ge \tilde{c} \, \nu_0^{\frac{3}{2p-3}} \, \nu_1^{\frac{-2}{2p-3}} \lambda_1^{\frac{1}{2p-3}} \, G^{\frac{4}{2p-3}}.$$

and h > 0 small enough such that

$$\mu c_0^2 h^2 \le \nu_0,$$

where \tilde{c} is a dimensionless number depending only on p and Ω , while c_0 is dimensionless constant given in (2.1), we have

$$\|\mathbf{u}(t) - \mathbf{v}(t)\| \le \|\mathbf{u}(t_0) - \mathbf{v}(t_0)\| e^{(\beta^* + M_0)} e^{-\nu_0 \lambda_1 \beta^* (t - t_0)}, \quad \forall t \ge t_0,$$

where the positive parameters β^* and M_0 are defined in (4.19) and (4.21), respectively, and t_0 is given in Proposition 3.3.

Proof. Subtracting (3.1) and (4.1), the difference e = u - v satisfies the following error equation

$$\langle \partial_t \mathbf{e}, \mathbf{w} \rangle + ((\mathbf{u} \cdot \nabla)\mathbf{u}, \mathbf{w}) - ((\mathbf{v} \cdot \nabla)\mathbf{v}, \mathbf{w}) + (\mathbf{T}(\mathbf{D}\mathbf{u}) - \mathbf{T}(\mathbf{D}\mathbf{v}), \mathbf{D}\mathbf{w}) = -\mu (I_h \mathbf{e}, \mathbf{w}). \tag{4.12}$$

Since

$$(\mathbf{u}\cdot\nabla)\mathbf{u}-(\mathbf{v}\cdot\nabla)\mathbf{v}=(\mathbf{e}\cdot\nabla)\mathbf{u}+(\mathbf{v}\cdot\nabla)\mathbf{e},$$

taking $\mathbf{w} = \mathbf{e}$ and using [38, Lemma 2.45], the Korn inequality and (3.3), we obtain

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{dt}} \|\mathbf{e}\|^2 + \nu_0 \|\nabla \mathbf{e}\|^2 \le -((\mathbf{e} \cdot \nabla)\mathbf{u}, \, \mathbf{e}) - (\mu I_h \mathbf{e}, \, \mathbf{e}). \tag{4.13}$$

In light of (2.1) and the assumption $\mu c_0^2 h^2 \le \nu_0$, one can estimate the nudging term in (4.13) as

$$-\mu (I_{h}\mathbf{e}, \mathbf{e}) = -\mu (I_{h}\mathbf{e} - \mathbf{e} + \mathbf{e}, \mathbf{e})$$

$$= \mu (\mathbf{e} - I_{h}\mathbf{e}, \mathbf{e}) - \mu \|\mathbf{e}\|^{2}$$

$$\leq \frac{\mu}{2} \|\mathbf{e} - I_{h}\mathbf{e}\|^{2} + \frac{\mu}{2} \|\mathbf{e}\|^{2} - \mu \|\mathbf{e}\|^{2}$$

$$\leq \frac{\mu}{2} c_{0}^{2} h^{2} \|\nabla \mathbf{e}\|^{2} - \frac{\mu}{2} \|\mathbf{e}\|^{2}$$

$$\leq \frac{\nu_{0}}{2} \|\nabla \mathbf{e}\|^{2} - \frac{\mu}{2} \|\mathbf{e}\|^{2}$$

$$\leq \frac{\nu_{0}}{2} \|\nabla \mathbf{e}\|^{2} - \frac{\mu}{2} \|\mathbf{e}\|^{2}.$$
(4.14)

Thus, we have

$$\frac{1}{2} \frac{d}{dt} \|\mathbf{e}\|^2 + \frac{\nu_0}{2} \|\nabla \mathbf{e}\|^2 \le |((\mathbf{e} \cdot \nabla) \mathbf{u}, \mathbf{e})| - \frac{\mu}{2} \|\mathbf{e}\|^2.$$
 (4.15)

Take p and p' to be conjugate numbers, i.e., $p' = \frac{p}{p-1}$, and apply the (Lebesgue interpolation inequality), (Sobolev embedding) and (Young inequality) to estimate the above nonlinear term as

$$| ((\mathbf{e} \cdot \nabla) \mathbf{u}, \mathbf{e}) | \leq \| \mathbf{e}^{2} \|_{L^{p'}} \| \nabla \mathbf{u} \|_{L^{p}} = \| \mathbf{e} \|_{L^{2p'}}^{2} \| \nabla \mathbf{u} \|_{L^{p}} \leq \| \mathbf{e} \|^{2 - \frac{3}{p}} \| \mathbf{e} \|_{L^{6}}^{\frac{3}{p}} \| \nabla \mathbf{u} \|_{L^{p}}$$

$$\leq c_{S}^{\frac{3}{p}} \| \mathbf{e} \|^{2 - \frac{3}{p}} \| \nabla \mathbf{e} \|^{\frac{3}{p}} \| \nabla \mathbf{u} \|_{L^{p}} \leq \frac{\nu_{0}}{2} \| \nabla \mathbf{e} \|^{2} + \frac{\bar{c}}{2} \nu_{0}^{\frac{3}{3 - 2p}} \| \nabla \mathbf{u} \|_{L^{p}}^{\frac{2p}{2p - 3}} \| \mathbf{e} \|^{2},$$

$$(4.16)$$

for some \bar{c} depending only on p and Ω . Inserting (4.16) in (4.15), we get

$$\frac{\mathrm{d}}{\mathrm{d}t} \|\mathbf{e}\|^2 + \left(\mu - \bar{c}\,\nu_0^{\frac{3}{3-2p}} \|\nabla\mathbf{u}\|_{L^p}^{\frac{2p}{2p-3}}\right) \|\mathbf{e}\|^2 \le 0. \tag{4.17}$$

With Lemma 2.5 in mind, denote

$$\alpha(t) = \mu - \bar{c} \nu_0^{\frac{3}{3-2p}} \|\nabla \mathbf{u}(t)\|_{L^p}^{\frac{2p}{2p-3}}.$$

Applying Hölder's inequality, and choosing $T = (\nu_0 \lambda_1)^{-1}$ in (3.8), we obtain for $p \geq \frac{5}{2}$

$$\begin{split} \int_{t}^{t+T} \alpha(s) \, ds &= \mu T - \bar{c} \, \nu_{0}^{\frac{3}{3-2p}} \int_{t}^{t+T} \| \nabla \mathbf{u}(s) \|_{L^{p}}^{\frac{2p}{2p-3}} ds \\ &\geq \mu T - \bar{c} \, \nu_{0}^{\frac{3}{3-2p}} \, T^{\frac{2p-5}{2p-3}} \left(\int_{t}^{t+T} \| \nabla \mathbf{u}(s) \|_{L^{p}}^{p} \, ds \right)^{\frac{2}{2p-3}} \\ &\geq \mu T - \bar{c} \, \nu_{0}^{\frac{3}{3-2p}} \, T^{\frac{2p-5}{2p-3}} \left(2c_{K}^{p} \, (1 + \nu_{0} \, \lambda_{1} \, T) \, \nu_{0}^{2} \, \nu_{1}^{-1} \, \lambda_{1}^{-\frac{1}{2}} \, G^{2} \right)^{\frac{2}{2p-3}} \\ &= \frac{\mu}{\nu_{0} \, \lambda_{1}} - 2^{\frac{4}{2p-3}} \bar{c} \, c_{K}^{\frac{2p}{2p-3}} \, \nu_{0}^{\frac{6-2p}{2p-3}} \, \nu_{1}^{\frac{-2}{2p-3}} \, \lambda_{1}^{\frac{4-2p}{2p-3}} \, G^{\frac{4}{2p-3}}. \end{split}$$

Thus, from above and with $\mu \geq 2^{1+\frac{4}{2p-3}} \bar{c} c_K^{\frac{2p}{2p-3}} \nu_0^{\frac{3}{2p-3}} \nu_1^{\frac{-2}{2p-3}} \lambda_1^{\frac{1}{2p-3}} G^{\frac{4}{2p-3}}$, we have

$$\int_{t}^{t+T} \alpha(s) \, ds \ge \beta^* > 0, \qquad \forall \, t \ge t_0, \tag{4.18}$$

where

$$\beta^* := 2^{\frac{4}{2p-3}} \bar{c} c_K^{\frac{2p}{2p-3}} \nu_0^{\frac{6-2p}{2p-3}} \nu_1^{\frac{-2}{2p-3}} \left(\frac{1}{\lambda_1}\right)^{\frac{2p-4}{2p-3}} G^{\frac{4}{2p-3}}. \tag{4.19}$$

Similarly, by setting $\alpha^{-}(s) = \max\{-\alpha(s), 0\}$, we have

$$\int_{t}^{t+T} \alpha^{-}(s) ds \le M_0, \qquad \forall t \ge t_0, \tag{4.20}$$

where

$$M_0 := 2^{\frac{4}{2p-3}} \bar{c} \, c_K^{\frac{2p}{2p-3}} \, \nu_0^{\frac{6-2p}{2p-3}} \, \nu_1^{\frac{-2}{2p-3}} \, \lambda_1^{\frac{4-2p}{2p-3}} \, G^{\frac{4}{2p-3}}. \tag{4.21}$$

Finally, by applying Lemma 2.5 to (4.17), we conclude that

$$\|\mathbf{e}(t)\| = \|\mathbf{u}(t) - \mathbf{v}(t)\| \le \|\mathbf{u}(t_0) - \mathbf{v}(t_0)\| e^{(\beta^* + M_0)} e^{-(\frac{t - t_0}{T})\beta^*}, \quad \forall t \ge t_0.$$

That is, the error converges exponentially fast to 0 as $t \to \infty$. \square

5. The case $p = \frac{11}{5}$ with periodic boundary conditions

In this section we study the dynamics of the solutions \mathbf{u} for the Ladyzhenskaya model $(3.1)_{1-2}$ and \mathbf{v} for the corresponding data assimilation algorithm $(4.1)_{1-2}$ in $\Omega = [0, 2\pi]^3$ completed with periodic boundary conditions.

Since the average velocity $\overline{\mathbf{u}}(t) = \int_{\Omega} \mathbf{u}(\mathbf{x}, t) d\mathbf{x}$ is an invariant of the flow provided that $\int_{\Omega} \mathbf{f}(\mathbf{x}, t) d\mathbf{x} = 0$ and the interpolant operators (volume elements or Fourier modes) have zero spatial average, we consider without loss of generality that $\overline{\mathbf{u}}(t) = 0$ and $\overline{\mathbf{v}}(t) = 0$ for all $t \geq 0$.

Theorem 5.1 (Existence of Weak Solutions and their Propagation of Regularity). Let $p = \frac{11}{5}$, $\mathbf{f} \in L^2(0,T;\dot{L}^2(\Omega))$ and $\mathbf{u}_0 \in \dot{L}^2_{\sigma}(\Omega)$. Then, there exists a weak solution \mathbf{u} to $(3.1)_{1-2}$ on $(0,\infty)$ with periodic boundary conditions such that

$$\mathbf{u} \in \mathcal{C}([0,T]; \dot{L}_{\sigma}^{2}(\Omega)) \cap L^{\frac{11}{5}}(0,T; W_{\sigma}^{1,\frac{11}{5}}(\Omega)), \quad \partial_{t}\mathbf{u} \in L^{\frac{11}{6}}(0,T; (W_{\sigma}^{1,\frac{11}{5}}(\Omega))'), \qquad \forall T \ge 0,$$
 (5.1)

and

$$\langle \partial_t \mathbf{u}, \mathbf{w} \rangle + ((\mathbf{u} \cdot \nabla)\mathbf{u}, \mathbf{w}) + (\mathbf{T}(\mathbf{D}\mathbf{u}), \nabla \mathbf{w}) = (\mathbf{f}, \mathbf{w}), \quad \forall \, \mathbf{w} \in W_{\sigma}^{1, \frac{11}{5}}(\Omega),$$
 (5.2)

for almost all $t \in [0,T]$. Moreover, the energy equality holds

$$\frac{1}{2} \|\mathbf{u}(t)\|^{2} + \int_{0}^{t} \left(2\nu_{0} \|\mathbf{D}\mathbf{u}(\tau)\|^{2} + 2\nu_{1} \|\mathbf{D}\mathbf{u}(\tau)\|^{\frac{11}{5}}_{L^{\frac{11}{5}}} \right) d\tau = \frac{1}{2} \|\mathbf{u}_{0}\|^{2} + \int_{0}^{t} (\mathbf{f}(\tau), \mathbf{u}(\tau)) d\tau, \qquad \forall t \ge 0. \quad (5.3)$$

In particular, if $\mathbf{f} \in \dot{L}^2(\Omega)$, there exists a time $t_0 > 0$ such that for all $t \geq t_0$ we have

$$\|\mathbf{u}(t)\|^2 \le 2\frac{\nu_0^2 G^2}{\lambda_1^{\frac{1}{2}}} \tag{5.4}$$

and

$$\int_{t}^{t+T} \left(\nu_0 \| \mathbf{D} \mathbf{u}(\tau) \|^2 + \nu_1 \| \mathbf{D} \mathbf{u}(\tau) \|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \right) d\tau \le 2 \left(1 + \nu_0 \lambda_1 T \right) \frac{\nu_0^2 G^2}{\lambda_1^{\frac{1}{2}}}, \tag{5.5}$$

where G is defined as in (3.5). In addition, there exists $\bar{t} \in [t_0, t_0 + 1]$ such that

$$\mathbf{u} \in L^{\infty}(\bar{t}, T; \dot{H}^{1}_{\sigma}(\Omega)) \cap L^{2}(\bar{t}, T; H^{2}_{\sigma}(\Omega)) \cap L^{\frac{11}{5}}(\bar{t}, T; W^{1, \frac{33}{5}}(\Omega)), \qquad \forall T \ge \bar{t}, \tag{5.6}$$

and

$$\int_{t}^{t+r} \|\nabla \mathbf{u}(\tau)\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} d\tau \le \frac{1}{K_{1}} \left(R_{3} + K_{2}R_{2}R_{3} + K_{3}R_{2} + \nu_{0}^{2} \lambda_{1}^{\frac{1}{2}} G^{2} \right), \qquad \forall t \ge t_{1}, \tag{5.7}$$

where $r = (\nu_0 \lambda_1)^{-1}$, $t_1 = \bar{t} + r$. The constants K_1 , K_2 , K_3 are defined in (5.19), and R_1 , R_2 , R_3 are given in (5.21)–(5.22).

Proof. The first part of Theorem 5.1 is proved in [38, Section 5] (see also [39] Theorem 3.1). Let us now consider a generic² weak solution \mathbf{u} to $(3.1)_{1-2}$ on $(0,\infty)$ satisfying (5.1), (5.2), (5.3), (5.4) and (5.5). It follows from (5.5) that there exists $\bar{t} \in [t_0, t_0 + 1]$ such that

$$\|\mathbf{D}\mathbf{u}(\bar{t})\| \le \left(2(1+\nu_0\lambda_1)\frac{\nu_0G^2}{\lambda_1^{\frac{1}{2}}}\right)^{\frac{1}{2}}.$$

Since $\mathbf{u}(\bar{t}) \in \dot{H}^{1}_{\sigma}(\Omega)$, we infer from [38, Theorem 3.4, Theorem 4.5 and Remark 4.6] (see also [39, Theorem 4.1]) that there exists a unique strong solution $\tilde{\mathbf{u}}$ on $[\bar{t}, \infty)$ originating from \mathbf{u} such that

$$\widetilde{\mathbf{u}} \in L^{\infty}(\overline{t}, T; \dot{H}^{1}_{\sigma}(\Omega)) \cap L^{2}(\overline{t}, T; H^{2}_{\sigma}(\Omega)) \cap L^{\frac{11}{5}}(\overline{t}, T; W^{1, \frac{33}{5}}(\Omega)), \qquad \forall T \geq \overline{t},$$

In addition, in light of the weak–strong uniqueness principle proved in [39, Theorem 5.2], we infer that $\tilde{\mathbf{u}}(t) = \mathbf{u}(t)$ for any $t \in [\bar{t}, \infty)$. This, in turn, gives (5.6).

We now perform some formal Sobolev estimates whose rigorous justification can be performed through the Galerkin scheme. By definition of the Stokes operator in the periodic setting, multiplying $(4.1)_1$ by $-\Delta \mathbf{u}$ and integrating over Ω , we obtain

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{dt}} \|\nabla \mathbf{u}\|^{2} + \nu_{0} \|\Delta \mathbf{u}\|^{2} + 2\nu_{1} \int_{\Omega} \nabla \cdot (|\mathrm{D}\mathbf{u}|_{F}^{\frac{1}{5}} \mathrm{D}\mathbf{u}) \cdot \Delta \mathbf{u} \, d\mathbf{x}
= -\int_{\Omega} \mathbf{f} \cdot \Delta \mathbf{u} \, d\mathbf{x} + \int_{\Omega} (\mathbf{u} \cdot \nabla) \mathbf{u} \cdot \Delta \mathbf{u} \, d\mathbf{x}.$$
(5.8)

Here we have used that $\nabla \cdot ((\nabla \mathbf{u})^T) = \nabla (\nabla \cdot \mathbf{u}) = 0$ by $(4.1)_2$. A direct calculation shows that

$$\partial_k(|\mathbf{D}\mathbf{u}|_F^n) = n|\mathbf{D}\mathbf{u}|_F^{n-2}\mathbf{D}\mathbf{u} : \mathbf{D}(\partial_k\mathbf{u}), \quad \forall n > 0.$$
 (5.9)

Using integration by parts and (5.9) with n = p - 2, we have for $p \ge 2$

$$\int_{\Omega} \nabla \cdot \left(|\operatorname{D}\mathbf{u}|_{F}^{p-2} \operatorname{D}\mathbf{u} \right) \cdot \Delta \mathbf{u} \, d\mathbf{x} = \int_{\Omega} \partial_{j} \left(|\operatorname{D}\mathbf{u}|_{F}^{p-2} (\operatorname{D}\mathbf{u})_{ij} \right) \partial_{kk} \mathbf{u}_{i} \, d\mathbf{x}
= -\int_{\Omega} |\operatorname{D}\mathbf{u}|_{F}^{p-2} (\operatorname{D}\mathbf{u})_{ij} \partial_{kk} \partial_{j} \mathbf{u}_{i} \, d\mathbf{x}
= \int_{\Omega} \partial_{k} \left(|\operatorname{D}\mathbf{u}|_{F}^{p-2} (\operatorname{D}\mathbf{u})_{ij} \right) \partial_{k} (\operatorname{D}\mathbf{u})_{ij} \, d\mathbf{x}
= \int_{\Omega} \partial_{k} \left(|\operatorname{D}\mathbf{u}|_{F}^{p-2} (\operatorname{D}\mathbf{u})_{ij} \partial_{k} (\operatorname{D}\mathbf{u})_{ij} \, d\mathbf{x} + \int_{\Omega} |\operatorname{D}\mathbf{u}|_{F}^{p-2} \partial_{k} (\operatorname{D}\mathbf{u})_{ij} \partial_{k} (\operatorname{D}\mathbf{u})_{ij} \, d\mathbf{x}
= \int_{\Omega} (p-2) |\operatorname{D}\mathbf{u}|_{F}^{p-4} (\operatorname{D}\mathbf{u})_{lm} (\operatorname{D}\partial_{k}\mathbf{u})_{lm} (\operatorname{D}\mathbf{u})_{ij} (\operatorname{D}\partial_{k}\mathbf{u})_{ij} \, d\mathbf{x} + \int_{\Omega} |\operatorname{D}\mathbf{u}|_{F}^{p-2} |\nabla(\operatorname{D}\mathbf{u})|^{2} \, d\mathbf{x}
= \int_{\Omega} (p-2) |\operatorname{D}\mathbf{u}|_{F}^{p-4} |\operatorname{D}\mathbf{u} : \operatorname{D}(\nabla\mathbf{u})|^{2} \, d\mathbf{x} + \int_{\Omega} |\operatorname{D}\mathbf{u}|_{F}^{p-2} |\nabla(\operatorname{D}\mathbf{u})|^{2} \, d\mathbf{x}. \tag{5.10}$$

Exploiting again (5.9) with $n = \frac{p}{2}$, we observe that

$$\int_{\mathcal{Q}} |\nabla |\mathrm{D}\mathbf{u}|_F^{\frac{p}{2}}|^2 d\mathbf{x} = \left(\frac{p}{2}\right)^2 \int_{\mathcal{Q}} |\mathrm{D}\mathbf{u}|_F^{p-4} |\mathrm{D}\mathbf{u} : \mathrm{D}(\nabla \mathbf{u})|^2 d\mathbf{x}.$$

As a consequence, it follows for $p = \frac{11}{5}$ that

$$\int_{\Omega} \nabla \cdot \left(|\mathrm{D}\mathbf{u}|_{F}^{\frac{1}{5}} \mathrm{D}\mathbf{u} \right) \cdot \Delta \mathbf{u} \, d\mathbf{x} \ge \frac{1}{5} \cdot \left(\frac{10}{11} \right)^{2} \int_{\Omega} |\nabla |\mathrm{D}\mathbf{u}|_{F}^{\frac{11}{10}}|^{2} \, d\mathbf{x}$$

² Indeed, in the case $p \in \left[\frac{11}{5}, \frac{5}{2}\right)$, the weak solutions are not known to be unique (cf. [39]).

$$\begin{split} &= \frac{20}{121} \left\| |\mathbf{D}\mathbf{u}|_F^{\frac{11}{10}} \right\|_{H^1}^2 - \frac{20}{121} \left\| |\mathbf{D}\mathbf{u}|_F^{\frac{11}{10}} \right\|^2 \\ &\geq \frac{1}{8} \left\| |\mathbf{D}\mathbf{u}|_F^{\frac{11}{10}} \right\|_{H^1}^2 - \frac{1}{6} \left\| \mathbf{D}\mathbf{u} \right\|_{L^{\frac{11}{5}}}^{\frac{11}{5}}. \end{split}$$

Using the embedding $H^1(\Omega) \hookrightarrow L^6(\Omega)$ and the Korn inequality, we infer that

$$\begin{split} \int_{\varOmega} \nabla \cdot \left(|\mathrm{D}\mathbf{u}|_{F}^{\frac{1}{5}} \mathrm{D}\mathbf{u} \right) \cdot \varDelta \mathbf{u} \, d\mathbf{x} &\geq \frac{1}{8} \frac{1}{c_{S}^{2}} \left\| |\mathrm{D}\mathbf{u}|_{F}^{\frac{11}{10}} \right\|_{L^{6}}^{2} - \frac{1}{6} \left\| \mathrm{D}\mathbf{u} \right\|_{L^{p}}^{p} \\ &\geq \frac{1}{8} \frac{C^{\frac{11}{5}}}{c_{S}^{2}} \left\| \mathrm{D}\mathbf{u} \right\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} - \frac{1}{6} \left\| \mathrm{D}\mathbf{u} \right\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \\ &\geq \frac{1}{8} \frac{C^{\frac{11}{5}}}{c_{S}^{2}} \left\| \nabla \mathbf{u} \right\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} - \frac{1}{6} \left\| \mathrm{D}\mathbf{u} \right\|_{L^{\frac{11}{5}}}^{\frac{11}{5}}. \end{split}$$

In order to handle the convective term, we observe that

$$\int_{\Omega} (\mathbf{u} \cdot \nabla) \mathbf{u} \cdot \Delta \mathbf{u} \, d\mathbf{x} = \int_{\Omega} \mathbf{u}_{j} \partial_{j} \mathbf{u}_{i} \partial_{k} \mathbf{u}_{i} \, d\mathbf{x}$$

$$= -\int_{\Omega} \partial_{k} \mathbf{u}_{j} \partial_{j} \mathbf{u}_{i} \partial_{k} \mathbf{u}_{i} \, d\mathbf{x} - \int_{\Omega} \mathbf{u}_{j} \partial_{j} \partial_{k} \mathbf{u}_{i} \partial_{k} \mathbf{u}_{i} \, d\mathbf{x}$$

$$= -\int_{\Omega} \partial_{k} \mathbf{u}_{j} \partial_{j} \mathbf{u}_{i} \partial_{k} \mathbf{u}_{i} \, d\mathbf{x} - \underbrace{\int_{\Omega} \mathbf{u}_{j} \partial_{j} \left(\frac{1}{2} \partial_{k} \mathbf{u}_{i} \partial_{k} \mathbf{u}_{i}\right) \, d\mathbf{x}}_{-0} \leq \|\nabla \mathbf{u}\|_{L^{3}}^{3}.$$
(5.11)

Thus, collecting the above terms together, we find the differential inequality

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{dt}} \|\nabla \mathbf{u}\|^2 + \nu_0 \|\Delta \mathbf{u}\|^2 + \frac{\nu_1 \widetilde{C}}{4} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} \le \|\nabla \mathbf{u}\|_{L^3}^{\frac{3}{5}} + \frac{\nu_1}{3} \|\mathrm{D}\mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} - \int_{\Omega} \mathbf{f} \cdot \Delta \mathbf{u} \, d\mathbf{x}. \tag{5.12}$$

Here, we have set $\widetilde{C} = \frac{C^{\frac{11}{5}}}{c_S^2 c_K^{\frac{1}{5}}}$, which depends only on Ω and the value $p = \frac{11}{5}$. We now proceed with the estimate of the terms on the right-hand side of (5.12). We exploit the splitting method devised in [38] for the L^3 -norm of $\nabla \mathbf{u}$ which follows from the Lebesgue interpolation. We recall that for $p \in [2,3]$

$$\|\nabla \mathbf{u}\|_{L^{3}} \leq \|\nabla \mathbf{u}\|_{L^{p}}^{\frac{p-1}{2}} \|\nabla \mathbf{u}\|_{L^{3p}}^{\frac{3-p}{2}}, \qquad \|\nabla \mathbf{u}\|_{L^{3}} \leq \|\nabla \mathbf{u}\|_{L^{2}}^{\frac{2p-2}{3p-2}} \|\nabla \mathbf{u}\|_{L^{3p}}^{\frac{p}{3p-2}}.$$

For $\alpha \in (0,1)$, which will be chosen later, exploiting the above interpolation inequalities, we obtain

$$\begin{split} \|\nabla \mathbf{u}\|_{L^{3}}^{3} &\leq \|\nabla \mathbf{u}\|_{L^{3}}^{3\alpha} \|\nabla \mathbf{u}\|_{L^{3}}^{3(1-\alpha)} \\ &\leq \|\nabla \mathbf{u}\|_{L^{p}}^{3\alpha^{\frac{p-1}{2}}} \|\nabla \mathbf{u}\|_{L^{3p}}^{3\alpha^{\frac{3-p}{2}}} \|\nabla \mathbf{u}\|_{L^{2}}^{3(1-\alpha)^{\frac{2p-2}{3p-2}}} \|\nabla \mathbf{u}\|_{L^{3p}}^{3(1-\alpha)^{\frac{p}{3p-2}}} \\ &\leq \|\nabla \mathbf{u}\|_{L^{p}}^{3\alpha^{\frac{p-1}{2}}} \|\nabla \mathbf{u}\|_{L^{2}}^{3(1-\alpha)^{\frac{2p-2}{3p-2}}} \|\nabla \mathbf{u}\|_{L^{3p}}^{3\alpha^{\frac{3-p}{2}}+3(1-\alpha)^{\frac{p}{3p-2}}}. \end{split}$$

$$(5.13)$$

In particular, for $p = \frac{11}{5}$, we have

$$\|\nabla \mathbf{u}\|_{L^{3}}^{3} \leq \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{9}{5}\alpha} \|\nabla \mathbf{u}\|_{L^{2}}^{\frac{36}{23}(1-\alpha)} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{33}{23}-\alpha\frac{27}{115}}.$$

Setting

$$\alpha = \frac{22}{45}, \quad s = \frac{5}{3}, \quad s' = \frac{5}{2},$$

and using the Young inequality, it follows that for any $\varepsilon > 0$

$$\|\nabla \mathbf{u}\|_{L^{3}}^{3} \leq \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{22}{25}} \|\nabla \mathbf{u}\|_{L^{2}}^{\frac{4}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{\frac{33}{25}}$$

$$\leq \frac{3\varepsilon}{5} \|\nabla \mathbf{u}\|_{L^{\frac{13}{5}}}^{\frac{13}{5}} + \frac{2}{5\varepsilon^{\frac{3}{2}}} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{2}.$$

$$(5.14)$$

Choosing $\varepsilon = \frac{5}{24}\nu_1\widetilde{C}$, we are led to

$$\|\nabla \mathbf{u}\|_{L^{3}}^{3} \leq \frac{\nu_{1}\widetilde{C}}{8} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} + \frac{2}{5} \left(\frac{24}{5\nu_{1}\widetilde{C}}\right)^{\frac{3}{2}} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{2}. \tag{5.15}$$

Also, we have

$$-\int_{\Omega} \mathbf{f} \cdot \Delta \mathbf{u} \, d\mathbf{x} \le \frac{\nu_0}{2} \|\Delta \mathbf{u}\|^2 + \frac{1}{2\nu_0} \|\mathbf{f}\|^2. \tag{5.16}$$

Combining (5.12) with (5.15) and (5.16), we end up with

$$\frac{1}{2} \frac{\mathrm{d}}{\mathrm{d}t} \|\nabla \mathbf{u}\|^{2} + \frac{\nu_{0}}{2} \|\Delta \mathbf{u}\|^{2} + \frac{\nu_{1} \widetilde{C}}{8} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} \\
\leq \frac{2}{5} \left(\frac{24}{5\nu_{1}\widetilde{C}}\right)^{\frac{3}{2}} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{2} + \frac{\nu_{1}}{3} \|\mathbf{D}\mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} + \frac{1}{2\nu_{0}} \|\mathbf{f}\|^{2}, \tag{5.17}$$

for almost any $t \in (\bar{t}, \infty)$. We rewrite the above inequality as

$$\frac{\mathrm{d}}{\mathrm{dt}} \|\nabla \mathbf{u}\|^{2} + \nu_{0} \|\Delta \mathbf{u}\|^{2} + K_{1} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} \le K_{2} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{2} + K_{3} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} + \frac{1}{\nu_{0}} \|\mathbf{f}\|^{2}, \tag{5.18}$$

having set

$$K_1 = \frac{\nu_1 \widetilde{C}}{4}, \quad K_2 = \frac{2}{5} \left(\frac{24}{5\nu_1 \widetilde{C}}\right)^{\frac{3}{2}}, \quad K_3 = \frac{2\nu_1 C}{3}.$$
 (5.19)

In particular, we have

$$\frac{\mathrm{d}}{\mathrm{d}\mathbf{t}} \|\nabla \mathbf{u}\|^{2} \le K_{2} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} \|\nabla \mathbf{u}\|_{L^{2}}^{2} + K_{3} \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} + \frac{1}{\nu_{0}} \|\mathbf{f}\|^{2}. \tag{5.20}$$

In light of (5.5), for any $t \ge t_0$ and $r = (\nu_0 \lambda_1)^{-1}$ we infer that

$$\int_{t}^{t+r} \|\nabla \mathbf{u}(\tau)\|^{2} d\tau \leq 8 \frac{\nu_{0} G^{2}}{\lambda_{1}^{\frac{1}{2}}} =: R_{1}, \qquad \int_{t}^{t+r} \|\nabla \mathbf{u}(\tau)\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} d\tau \leq 4 c_{K}^{\frac{11}{5}} \frac{\nu_{0}^{2} G^{2}}{\nu_{1} \lambda_{1}^{\frac{1}{2}}} =: R_{2}.$$
 (5.21)

By exploiting Lemma 2.4, we find

$$\|\nabla \mathbf{u}(t)\|^{2} \le \left(\nu_{0}\lambda_{1}R_{1} + K_{3}R_{2} + \nu_{0}^{2}\lambda_{1}^{\frac{1}{2}}G^{2}\right) e^{K_{2}R_{2}} =: R_{3}, \qquad \forall t \ge \bar{t} + r = t_{1}.$$
 (5.22)

As an immediate consequence, integrating (5.17) from t to t + r, where $t \ge t_1$, we obtain

$$\int_{t}^{t+r} \|\nabla \mathbf{u}(\tau)\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} d\tau \le \frac{1}{K_{1}} \left(R_{3} + K_{2}R_{2}R_{3} + K_{3}R_{2} + \nu_{0}^{2} \lambda_{1}^{\frac{1}{2}} G^{2} \right). \quad \Box$$
 (5.23)

Next, we state the following result concerning the existence of solutions to the data assimilation algorithm given by (4.1) in the case $p = \frac{11}{5}$. This is a consequence of the results obtained in [38,39].

Theorem 5.2 (Existence of Weak and Strong Solutions for Data Assimilation Problem). Assume that $p = \frac{11}{5}$ and $\mathbf{f} \in \dot{L}^2(\Omega)$. Let \mathbf{u} be a weak solution of (3.1) with periodic boundary conditions given by Theorem 5.1. Then, we have the following:

1. If $\mathbf{v}_0 \in \dot{L}^2_{\sigma}(\Omega)$, there exists a weak solution \mathbf{v} to (4.1) satisfying

$$\mathbf{v} \in \mathcal{C}([0,T]; \dot{L}_{\sigma}^{2}(\Omega)) \cap L^{\frac{11}{5}}(0,T; W_{\sigma}^{1,\frac{11}{5}}(\Omega)), \quad \partial_{t}\mathbf{v} \in L^{\frac{11}{6}}(0,T; (W_{\sigma}^{1,\frac{11}{5}}(\Omega))'), \qquad \forall T \ge 0, \quad (5.24)$$

and

$$\langle \partial_t \mathbf{v}, \mathbf{w} \rangle + ((\mathbf{v} \cdot \nabla) \mathbf{v}, \mathbf{w}) + (\mathbf{T}(\mathbf{D}\mathbf{v}), \nabla \mathbf{w}) = (\mathbf{f}, \mathbf{w}) - \mu(I_h(\mathbf{v} - \mathbf{u}), \mathbf{w}), \quad \forall \, \mathbf{w} \in \dot{W}_{\sigma}^{1, \frac{11}{5}}(\Omega), \tag{5.25}$$

for almost all $t \in [0, T]$.

2. If $\mathbf{v}_0 \in \dot{H}^1_{\sigma}(\Omega)$, there exists a unique strong solution \mathbf{v} to (4.1) such that

$$\mathbf{v} \in \mathcal{C}([0,T]; \dot{H}_{\sigma}^{1}(\Omega)) \cap L^{2}(0,T; H^{2}(\Omega)) \cap L^{\frac{11}{5}}(0,T; W^{1,\frac{33}{5}}(\Omega)), \qquad \forall T \ge 0, \tag{5.26}$$

which solves (4.1) in weak sense as in (5.25).

3. If $\mathbf{v}_0 \in \dot{W}_{\sigma}^{1,\frac{11}{5}}(\Omega)$, there exists a unique strong solution \mathbf{v} to (4.1) which satisfies, in addition to (5.26),

$$\mathbf{v} \in \mathcal{C}([0,T]; \dot{W}_{\sigma}^{1,\frac{11}{5}}(\Omega)), \quad \partial_t \mathbf{v} \in L^2(0,T; \dot{L}_{\sigma}^2(\Omega)), \qquad \forall T \ge 0.$$
 (5.27)

In particular, in this case \mathbf{v} solves (4.1) in weak sense with $\mathbf{w} \in \dot{H}^{1}_{\sigma}(\Omega)$.

Lastly, we prove the convergence result for $p = \frac{11}{5}$ in the periodic boundary setting.

Theorem 5.3. For $p = \frac{11}{5}$, let **u** be a weak solution of (3.1) with periodic boundary conditions given by Theorem 5.1 and let **v** be the solution to the data assimilation algorithm given by Theorem 5.2. Assume that

$$\mu \ge \frac{2\overline{C}\nu_0^{\frac{5}{17}}\lambda_1^{\frac{10}{17}}}{K_1^{\frac{10}{17}}} \left(R_3 + K_2R_2R_3 + K_3R_2 + \nu_0^2\lambda_1^{\frac{1}{2}}G^2 \right)^{\frac{10}{17}}$$
(5.28)

where \overline{C} is a constant depending on Ω and p and K_1, K_2, K_3, R_2, R_3 are defined in Theorem 5.1, and h small enough such that

$$\mu c_0 h^2 \le \nu_0,$$

where c_0 is a dimensionless constant given (2.1). Then, we have

$$\|\mathbf{u}(t) - \mathbf{v}(t)\| \le \|\mathbf{u}(t_1) - \mathbf{v}(t_1)\| e^{(\gamma^* + M_1)} e^{-\nu_0 \lambda_1 \gamma^* (t - t_1)}, \quad \forall t \ge t_1,$$

where t_1 is as in (5.22) and the positive parameters γ^* and M_1 are defined in (5.32) and in (5.33), respectively.

Proof. Proceeding as in the proof of Theorem 4.2, we have

$$\frac{1}{2}\frac{\mathrm{d}}{\mathrm{d}t}\|\mathbf{e}\|^2 + \frac{\nu_0}{2}\|\nabla\mathbf{e}\|^2 \le |((\mathbf{e}\cdot\nabla)\mathbf{u},\mathbf{e})| - \frac{\mu}{2}\|\mathbf{e}\|^2.$$
(5.29)

Arguing differently than (4.16), we find

$$\begin{split} |\left(\left(\mathbf{e}\cdot\nabla\right)\mathbf{u},\mathbf{e}\right)| &\leq \|\mathbf{e}^{2}\|_{L^{\frac{33}{28}}} \|\nabla\mathbf{u}\|_{L^{\frac{33}{5}}} = \|\mathbf{e}\|_{L^{\frac{33}{14}}}^{2} \|\nabla\mathbf{u}\|_{L^{\frac{33}{5}}} \leq \|\mathbf{e}\|_{L^{\frac{11}{11}}}^{\frac{17}{11}} \|\mathbf{e}\|_{L^{\frac{16}{15}}}^{\frac{5}{17}} \|\nabla\mathbf{u}\|_{L^{\frac{33}{5}}} \\ &\leq c_{S}^{\frac{5}{11}} \|\mathbf{e}\|_{L^{\frac{17}{11}}}^{\frac{17}{11}} \|\nabla\mathbf{e}\|_{L^{\frac{5}{11}}}^{\frac{5}{11}} \|\nabla\mathbf{u}\|_{L^{\frac{33}{5}}} \leq \frac{\nu_{0}}{2} \|\nabla\mathbf{e}\|^{2} + c_{S}^{\frac{10}{17}} \left(\frac{2}{\nu_{0}}\right)^{\frac{5}{17}} \|\nabla\mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{22}{17}} \|\mathbf{e}\|^{2}. \end{split} \tag{5.30}$$

Inserting (5.30) in (5.29), we arrive at

$$\frac{\mathrm{d}}{\mathrm{dt}} \|\mathbf{e}\|^2 + \left(\mu - \frac{\overline{C}}{\nu_0^{\frac{5}{17}}} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{22}{17}}\right) \|\mathbf{e}\|^2 \le 0, \tag{5.31}$$

for some constant \overline{C} depending only on Ω and the value $p=\frac{11}{5}$. Aiming to use Lemma 2.5, let us set

$$\alpha(t) = \left(\mu - \frac{\overline{C}}{\nu_0^{\frac{5}{17}}} \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{22}{17}}\right).$$

By Hölder's inequality and (5.23), and also taking $r = (\nu_0 \lambda_1)^{-1}$, we obtain

$$\begin{split} \int_{t}^{t+r} \alpha(s) \, ds &= \mu r - \frac{\overline{C}}{\nu_{0}^{\frac{5}{17}}} \int_{t}^{t+r} \| \nabla \mathbf{u}(s) \|_{L^{\frac{22}{13}}}^{\frac{22}{17}} \, ds \\ &\geq \frac{\mu}{\nu_{0} \lambda_{1}} - \frac{\overline{C}}{\nu_{0}^{\frac{5}{17}}} \left(\int_{t}^{t+r} \| \nabla \mathbf{u}(s) \|_{L^{\frac{33}{5}}}^{\frac{11}{5}} \, ds \right)^{\frac{10}{17}} \left(\int_{t}^{t+r} 1 \, ds \right)^{\frac{7}{17}} \\ &\geq \frac{\mu}{\nu_{0} \lambda_{1}} - \frac{\overline{C}}{\nu_{0}^{\frac{12}{17}} \lambda_{1}^{\frac{7}{17}}} \frac{1}{K_{1}^{\frac{10}{17}}} \left(R_{3} + K_{2} R_{2} R_{3} + K_{3} R_{2} + \nu_{0}^{2} \lambda_{1}^{\frac{1}{2}} G^{2} \right)^{\frac{10}{17}}. \end{split}$$

Notice that the second term on the right-hand side of the above inequality is independent of μ . In particular, in light of the assumption (5.28), we immediately deduce that

$$\int_{t}^{t+r} \alpha(s) \, ds \ge \gamma^{\star}, \qquad \forall \, t \ge t_1,$$

where

$$\gamma^{\star} = \frac{2\overline{C}}{\nu_0^{\frac{12}{17}} \lambda_1^{\frac{7}{17}} K_1^{\frac{10}{17}}} \left(R_3 + K_2 R_2 R_3 + K_3 R_2 + \nu_0^2 \lambda_1^{\frac{1}{2}} G^2 \right)^{\frac{10}{17}}. \tag{5.32}$$

In a similar way, we find

$$\int_{t}^{t+r} \alpha^{-}(s) \, ds \le M_1, \qquad \forall \, t \ge t_1,$$

where

$$M_{1} = \frac{\overline{C}}{\nu_{0}^{\frac{12}{17}} \lambda_{1}^{\frac{7}{17}}} \frac{1}{K_{1}^{\frac{10}{17}}} \left(R_{3} + K_{2} R_{2} R_{3} + K_{3} R_{2} + \nu_{0}^{2} \lambda_{1}^{\frac{1}{2}} G^{2} \right)^{\frac{10}{17}}.$$
 (5.33)

Therefore, we conclude from Lemma 2.5 that

$$\|\mathbf{e}(t)\| = \|\mathbf{u}(t) - \mathbf{v}(t)\| \le \|\mathbf{u}(t_1) - \mathbf{v}(t_1)\| e^{(\gamma^* + M_1)} e^{-(\frac{t - t_1}{r})\gamma^*}, \quad \forall t \ge t_1. \quad \Box$$

Remark (2D Case). The condition (5.28) for the nudging parameter μ can be enhanced in 2D. Indeed, recalling that $\int_{\Omega} (\mathbf{u} \cdot \nabla) \mathbf{u} \cdot \Delta \mathbf{u} \, d\mathbf{x} = 0$, (5.18) is replaced by

$$\frac{\mathrm{d}}{\mathrm{dt}} \|\nabla \mathbf{u}\|^2 + \nu_0 \|\Delta \mathbf{u}\|^2 + 2K_1 \|\nabla \mathbf{u}\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} \le K_3 \|\nabla \mathbf{u}\|_{L^{\frac{11}{5}}}^{\frac{11}{5}} + \frac{1}{\nu_0} \|\mathbf{f}\|^2.$$
 (5.34)

Then, arguing as in the proof of Theorem 5.1, it follows that

$$\|\nabla \mathbf{u}(t)\|^{2} \le \left(\nu_{0}\lambda_{1}R_{1} + K_{3}R_{2} + \nu_{0}^{2}\lambda_{1}^{\frac{1}{2}}G^{2}\right) =: R_{3}^{\star}, \qquad \forall t \ge \bar{t} + r = t_{1},$$

$$(5.35)$$

and

$$\int_{t}^{t+r} \|\nabla \mathbf{u}(\tau)\|_{L^{\frac{33}{5}}}^{\frac{11}{5}} d\tau \le \frac{1}{2K_{1}} \left(R_{3}^{\star} + K_{3}R_{2} + \nu_{0}^{2} \lambda_{1}^{\frac{1}{2}} G^{2} \right), \qquad \forall t \ge t_{1}.$$
 (5.36)

As a direct consequence, (5.28) becomes

$$\mu \ge \frac{2\overline{C}\nu_0^{\frac{57}{17}}\lambda_1^{\frac{10}{17}}}{(2K_1)^{\frac{10}{17}}} \left(R_3^{\star} + K_3 R_2 + \nu_0^2 \lambda_1^{\frac{1}{2}} G^2 \right)^{\frac{10}{17}}.$$
 (5.37)

Furthermore, the analysis herein presented can be extended for any p > 2 in (1.4).

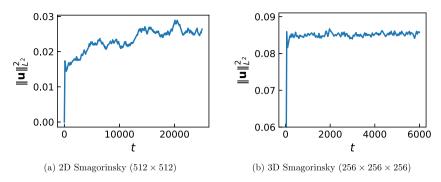


Fig. 1. Evolution of the energy of the reference solution over the transient period.

6. Computational results

We demonstrate the effectiveness of nudging for both two and three-dimensional Ladyzhenskaya models with fully periodic boundary conditions in $\Omega = [0, 2\pi]^d$, d = 2, 3. This is first done for the case p = 3, the Smagorinsky model, which is often used in Large Eddy Simulation (LES) of turbulent flow [46,47]. We then vary p in the three-dimensional case, and test nudging with only the horizontal components of velocity. For both cases, the parameter ν_1 is chosen from dimensional considerations to be

$$\nu_1 = \frac{1}{2} (C_s \delta)^2 \nu_0^{3-p} , \quad C_s = 0.1 , \quad \delta = \frac{2\pi}{N} ,$$
 (6.1)

where N is the number of Fourier modes used in each direction for the direct numerical simulation (DNS) of the reference solution.

The initial condition for the reference solution $\mathbf{u}(t_0)$ for each data assimilation experiment is chosen so that it faithfully reflects the long term dynamics of the model. This is done by integrating the model starting at t=0 with $\mathbf{u}(0)=0$ until some time $t=t_0$ when it appears the transient period has passed. Fig. 1 shows the time evolution of the energy $\|\mathbf{u}\|_{L^2}^2$ on $[0,t_0]$. By the end of the run, this quantity seems to have reached its statistically stationary state. We assume then that $\mathbf{u}(t_0)$ is essentially on the global attractor. We start the nudging at time $t=t_0$ by solving the original (\mathbf{u}) and the nudging (\mathbf{v}) systems simultaneously with $\mathbf{v}(t_0)=0$. The computations are done using Dedalus, an open-source spectral package (see [68]). The time stepper is a four-stage third order Runge–Kutta method.

6.1. Two-dimensional case

In two-dimensions, we take the viscosity to be $\nu_0=10^{-4},~\mu=1$, and use a normalized force $\mathbf{f}_{2\mathrm{D}}$ from [69], so that the Grashof number $G=2.5\times10^5$. We demonstrate both the nodal value and Fourier modes interpolant operators. In the nodal value case, we use every 4th nodal value in each direction so that $h\approx0.0491$. In the Fourier modes case, we use the projection on the low modes with wave vectors $\mathbf{k}=(k_1,k_2)$ such that $|k_j|\leq 32$ and $h=\frac{\pi}{32}\approx0.0982$. The value of N is fixed at 512. While we have not analyzed the nodal interpolation operator in this paper, Fig. 2(b) shows synchronization with the DNS of the reference solution to within machine precision in both the L^2 and H^1 norms. The same is true for Fourier mode interpolation, with a slower rate due to a larger value of h. Field plots of the velocity components and pressures at several times near the start of nudging corresponding to Fig. 2(a) are shown in Fig. 4.

6.2. Three-dimensional case

In the three-dimensional case, we define a force $\mathbf{f}_{3D} = (f_1, f_2, f_3)$ via its Fourier coefficients so that in each wave vector plane, \mathbf{f}_{3D} is similar to \mathbf{f}_{2D} in the previous section. Specifically, we take the function $g := \nabla \times \mathbf{f}_{2D}$

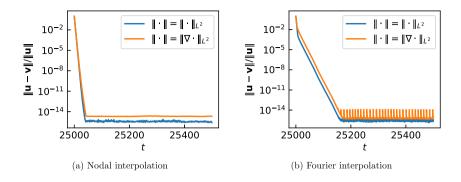


Fig. 2. Convergence of data assimilation for 2D Smagorinsky model for $\mu = 1$.

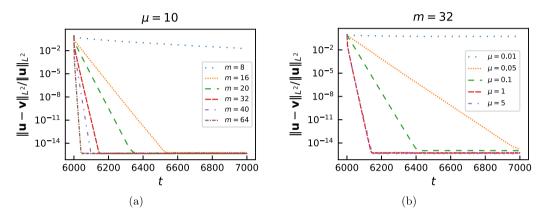


Fig. 3. Convergence of data assimilation for the 3D Smagorinsky at different values of the nudging parameter μ and h = h(m); the left fixes $\mu = 10$ and the right fixes m = 32.

and set

$$\begin{split} \hat{f}_1(k_1,0,k_3) &= \frac{ik_3\hat{g}(k_1,k_3)}{k_1^2 + k_3^2} \,, & \qquad \qquad \hat{f}_1(k_1,k_2,0) &= \frac{ik_2\hat{g}(k_1,k_2)}{k_1^2 + k_2^2} \,, \\ \hat{f}_2(k_1,k_2,0) &= \frac{-ik_1\hat{g}(k_1,k_2)}{k_1^2 + k_2^2} \,, & \qquad \qquad \hat{f}_2(0,k_2,k_3) &= \frac{ik_3\hat{g}(k_2,k_3)}{k_2^2 + k_3^2} \,, \\ \hat{f}_3(k_1,0,k_3) &= \frac{-ik_1\hat{g}(k_1,k_3)}{k_1^2 + k_3^2} \,, & \qquad \qquad \hat{f}_3(0,k_2,k_3) &= \frac{-ik_2\hat{g}(k_1,k_2)}{k_2^2 + k_3^2} \,, \end{split}$$

and all other Fourier coefficients of \mathbf{f}_{3D} are zero. In 3D it is the viscosity ν_0 that is adjusted so that the Grashof number remains as $G = 2.5 \times 10^5$. We use the Fourier modes interpolation operator $I_h = P_{h(m)}$ for the 3D model, where $P_{h(m)}$ denotes the projection on the low modes with wave vectors $\mathbf{k} = (k_1, k_2, k_3)$ such that $|k_j| \leq m$ and

$$h(m) = \frac{\pi}{m} \ .$$

The value of N is fixed at 256.

Fig. 3 shows the exponential rate of synchronization using different values of nudging parameter μ and resolution h. For fixed $\mu=10$, as we use fewer number of modes, the convergence is slower, but still exponential. For m=8, slices of solutions at the mid-plane $z=\pi$ near the start of nudging are shown in Fig. 5. The convergence fails at m=4 (not shown).

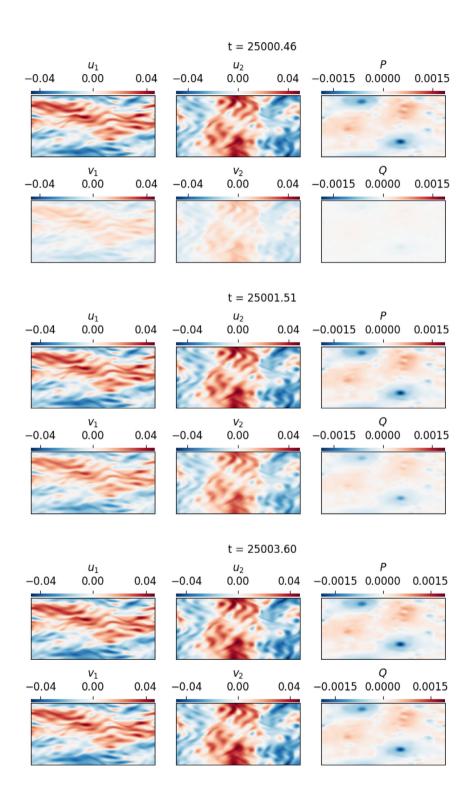


Fig. 4. Synchronization of the 2D Smagorinsky model using nodal interpolation, $\mu=1$ and $h\approx 0.0491$; the reference solution (\mathbf{u},P) is denoted as (u_1,u_2,P) and the nudging solution (\mathbf{v},Q) is (v_1,v_2,Q) .

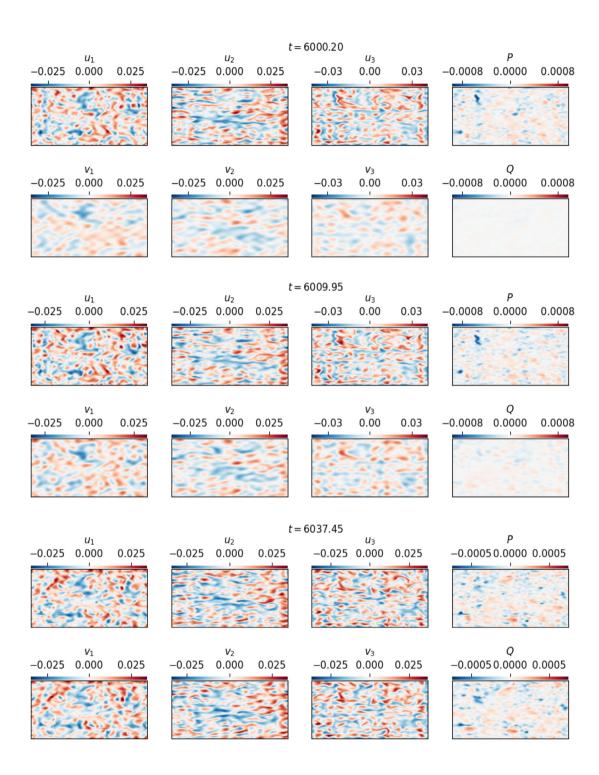


Fig. 5. Synchronization of the 3D Smagorinsky model using $\mu = 10, h = h(8)$. These are the slices in the mid-plane $(0, 2\pi) \times (0, 2\pi) \times \{z = \pi\}$.

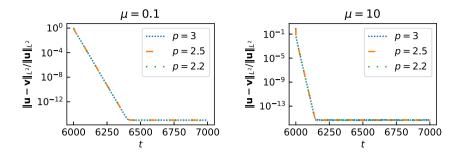


Fig. 6. Synchronization for the 3D Ladyzhenskaya model using h = h(32) for different values of p.

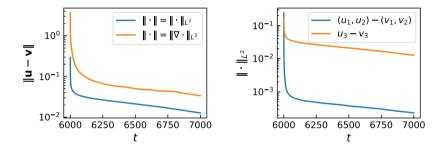


Fig. 7. Abridged nudging for the 3D Smagorinsky model with $(\mu_1, \mu_2, \mu_3) = (10, 10, 0)$ and h = h(128).

As is the case in other studies (e.g. [55,57,58]), the rigorous bounds on the parameters are not expected to be sharp, and in simulations μ can be taken much smaller and consequently the data much coarser than what is suggested by the analysis. At the fixed parameter of m=32, the convergence rate improves as μ is increased (see Fig. 3(b)). At $\mu=1$ and $\mu=5$, the convergence rates are nearly identical, while at $\mu=0.01$, nudging fails to synchronize. This numerical experiment suggests a critical value of μ .

We varied p (along with ν_1 according to (6.1)) in the Ladyzhenskaya model using both $\mu = 10$ and $\mu = 0.1$ (see Fig. 6). At these values of μ , we detect no discernible difference in the performance of the nudging algorithm for p ranging from 2.2 = 11/5 to 3.

Finally, we consider an abridged nudging scheme in which only the horizontal components of velocity play the role of observed data. This amounts to treating μ as the vector $(\mu_1, \mu_2, \mu_3) = (10, 10, 0)$ and nudging the jth component of velocity with the factor μ_j . Fig. 7 shows rapid initial synchronization, which then slows, particularly for the third component of velocity, which is not nudged. While the error is far from machine precision even after nudging for 1000 time units, the field plots shown in Fig. 8 display similar features at rates that are slower for the third component of velocity and pressure.

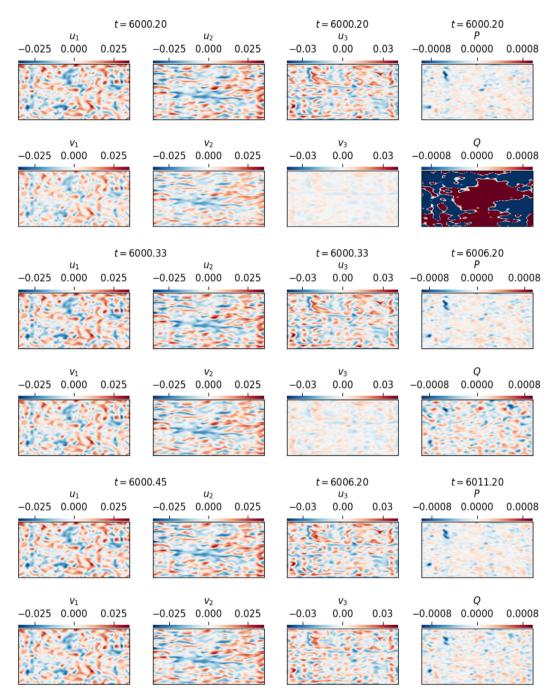


Fig. 8. Abridged data assimilation for the 3D Smagorinsky model with $(\mu_1, \mu_2, \mu_3) = (10, 10, 0)$ and h = h(128). These are the slices on the mid-plane $(0, 2\pi) \times (0, 2\pi) \times \{z = \pi\}$. Note the time progression is different for different components of velocity and pressure.

Declaration of competing interest

The authors declare that they have no known competing financial interests or personal relationships that could have appeared to influence the work reported in this paper.

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