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INCOMPRESSIBLE LIMIT FOR COMPRESSIBLE FLUIDS WITH STOCHASTIC FORCING

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ABSTRACT. We study the asymptotic behavior of the isentropic Navier–Stokes system driven by a multiplicative stochastic forcing in the compressible regime, where the Mach number approaches zero. Our approach is based on the recently developed concept of weak martingale solution to the primitive system, uniform bounds derived from a stochastic analogue of the modulated energy inequality, and careful analysis of acoustic waves. A stochastic incompressible Navier–Stokes system is identified as the limit problem.

1. INTRODUCTION

Singular limit processes bridge the gap between fluid motion considered in different geometries, times scales, and/or under different constitutive relations as the case may be. In their pioneering paper, Klainerman and Majda [15] proposed a general approach to these problems in the context of hyperbolic conservation laws, in particular, they examine the passage from compressible to incompressible fluid flow motion via the low Mach number limit. As the problems are typically non-linear, the method applies in general only on short time intervals on which regular solutions are known to exist. A qualitatively new way, at least in the framework of viscous fluids, has been open by the mathematical theory of weak solutions developed by P.-L. Lions [16]. In a series of papers, Lions and Masmoudi [17], [18] (see also Desjardins, Grenier [8], Desjardins et al. [9]) studied various singular limits for the *barotropic* Navier–Stokes system, among which the incompressible (low Mach number) limit. The incompressible limit is characterized with a large speed of the acoustic waves becoming infinite in the asymptotic regime. Accordingly, the fluid density approaches a constant and the velocity solenoidal. The limit behavior is described by the standard *incompressible* Navier–Stokes system.

In the present paper, we study the *compressible-incompressible* scenario in the context of stochastically driven fluids. Specifically, we consider the Navier–Stokes system for an isentropic compressible viscous fluid driven by a multiplicative stochastic forcing and study the asymptotic behavior of solutions in the low Mach number regime. To avoid the well known difficulties due to the presence of a boundary layer in the case of no-slip boundary conditions (cf. Desjardins et al. [9]), we restrict ourselves to the motion in the “flat” N -dimensional torus $\mathbb{T}^N = ([0, 2\pi]_{\{0, 2\pi\}})^N$, $N = 2, 3$ and on a finite time interval $(0, T)$; we set $Q = (0, T) \times \mathbb{T}^N$. We study the limit as $\varepsilon \rightarrow 0$ in the following system which governs the time evolution of the density ϱ and the velocity \mathbf{u} of a compressible viscous fluid:

$$(1.1a) \quad d\varrho + \operatorname{div}(\varrho \mathbf{u}) dt = 0,$$

$$(1.1b) \quad d(\varrho \mathbf{u}) + \left[\operatorname{div}(\varrho \mathbf{u} \otimes \mathbf{u}) - \nu \Delta \mathbf{u} - (\lambda + \nu) \nabla \operatorname{div} \mathbf{u} + \frac{1}{\varepsilon^2} \nabla p(\varrho) \right] dt = \Phi(\varrho, \varrho \mathbf{u}) dW.$$

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Here $p(\varrho)$ is the pressure which is supposed to follow the γ -law, i.e. $p(\varrho) = \varrho^\gamma$ where $\gamma > N/2$; the viscosity coefficients ν, λ satisfy

$$\nu > 0, \quad \lambda + \frac{2}{3}\nu \geq 0.$$

The driving process W is a cylindrical Wiener process defined on some probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and the coefficient Φ is a linear function of momentum $\varrho \mathbf{u}$ and a generally non-linear function of density ϱ satisfying suitable growth conditions. The precise description of the problem setting will be given in the next section.

In the limit we recover the stochastic Navier-Stokes system for incompressible fluids, that is,

$$(1.2a) \quad d\mathbf{u} + [\operatorname{div}(\mathbf{u} \otimes \mathbf{u}) - \nu \Delta \mathbf{u} + \nabla \pi] dt = \Psi(\mathbf{u}) dW,$$

$$(1.2b) \quad \operatorname{div}(\mathbf{u}) = 0,$$

where π denotes the associated pressure and $\Psi(\mathbf{u}) = \mathcal{P}_H \Phi(1, \mathbf{u})$, with \mathcal{P}_H being the Helmholtz projection onto the space of solenoidal vector fields. To be more precise, we show that for a given initial law Λ for (1.1) and the *ill-prepared* initial data for the compressible Navier-Stokes system (1.1), the approximate densities converge to a constant whereas the velocities converge *in law* to a weak martingale solution to the incompressible Navier-Stokes system (1.2) with the initial law Λ . This result is then strengthened in dimension two where we are able to prove the almost sure convergence of the velocities.

Our approach is based on the concept of *finite energy weak martingale solution* to the compressible Navier-Stokes system (1.1), whose existence was established recently in [2] and extends the approach in [11] to the stochastic setting, see Section 2 for more details. Similarly to its deterministic counterpart, the low Mach number limit problem features two essential difficulties:

- finding suitable uniform bounds independent of the scaling parameter ε ;
- analysis of rapidly oscillating *acoustic waves*, at least in the case of ill-prepared data.

Here, the necessary uniform bounds follow directly from the associated stochastic analogue of the energy inequality exploiting the basic properties of Itô's integral, see Section 3.1. The propagation of acoustic waves is described by a stochastic variant of Lighthill's acoustic analogy: A linear wave equation driven by a stochastic forcing, see Section 3.2. The desired estimates are obtained via the deterministic approach, specifically the so-called local method proposed by Lions and Masmoudi [17, 18], adapted to the stochastic setting.

A significant difference in comparison to the deterministic situation is the corresponding compactness argument. In general it is not possible to get any compactness in ω as no topological structure on the sample space Ω is assumed. To overcome this difficulty, it is classical to rather concentrate on compactness of the set of laws of the approximations and apply the Skorokhod representation theorem. It gives existence of a new probability space with a sequence of random variables that have the same laws as the original ones and that in addition converge almost surely. However, the Skorokhod representation Theorem is restricted to metric spaces but the structure of the compressible Navier-Stokes equations naturally leads to weakly converging sequences. On account of this we work with the Jakubowski-Skorokhod Theorem which is valid on a large class of topological spaces (including separable Banach spaces with weak topology). In the two-dimensional case we gain a stronger convergence result (see Theorem 2.9). This is based

on the uniqueness for the system (1.2) and a new version of the Gyöngy-Krylov characterization of convergence in probability [13] which applies to the setting of quasi-Polish spaces (see Proposition A.4).

We point out that the gradient part of the velocity converges only weakly to zero due to the presence of the acoustic waves, and, consequently, the limit in the stochastic forcing $\Phi(\varrho, \varrho \mathbf{u}) dW$ can be performed only if Φ is *linear* with respect to $\varrho \mathbf{u}$. However, this setting already covers the particular case of

$$\Phi(\varrho, \varrho \mathbf{u}) dW = \varrho \Phi_1 dW^1 + \varrho \mathbf{u} \Phi_2 dW^2$$

with two independent cylindrical Wiener processes W^1 and W^2 and suitable Hilbert-Schmidt operators Φ_1 and Φ_2 , which is the main example we have in mind. Here the first term describes some external force whereas the second one may be interpreted as a friction force of Brinkman's type, see e.g. Angot et al. [1].

In the case of $\Phi(\varrho, \varrho \mathbf{u}) = \varrho \Phi_1$, a semi-deterministic approach towards existence for (1.1) was developed in [10] (see also [23] for the two-dimensional case). More precisely, this particular case of multiplicative noise permits reduction of the problem that can be solved pathwise using deterministic arguments only. Nevertheless, it seems that such a pathwise approach is not convenient for the incompressible limit. In particular, uncontrolled quantities appear in the basic energy estimate and therefore the uniform bounds with respect to the parameter ε are lost. On the contrary, the stochastic method of the present paper heavily depends on the martingale properties of the Itô's stochastic integral which gives sufficient control of the expected values of all the necessary quantities.

The exposition is organized as follows. In Section 2 we continue with the introductory part: we introduce the basic set-up, the concept of solution and state the main results in Theorem 2.8 and Theorem 2.9. The remainder of the paper is then devoted to its proof.

2. MATHEMATICAL FRAMEWORK AND THE MAIN RESULT

Throughout the whole text, the symbols $W^{l,p}$ will denote the Sobolov space of functions having distributional derivatives up to order l integrable in L^p . We will also use $W^{l,2}(\mathbb{T}^N)$ for $l \in \mathbb{R}$ to denote the space of distributions v defined on \mathbb{T}^N with the finite norm

$$(2.1) \quad \sum_{k \in \mathbb{Z}} k^{2l} |c_k(v)|^2 < \infty,$$

where c_k denote the Fourier coefficients with respect to the standard trigonometric basis $\{\exp(ikx)\}_{k \in \mathbb{Z}}$.

To begin with, let us set up the precise conditions on the random perturbation of the system (1.1). Let $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t \geq 0}, \mathbb{P})$ be a stochastic basis with a complete, right-continuous filtration. The process W is a cylindrical Wiener process, that is, $W(t) = \sum_{k \geq 1} \beta_k(t) e_k$ with $(\beta_k)_{k \geq 1}$ being mutually independent real-valued standard Wiener processes relative to $(\mathcal{F}_t)_{t \geq 0}$ and $(e_k)_{k \geq 1}$ a complete orthonormal system in a separable Hilbert space \mathfrak{U} . To give the precise definition of the diffusion coefficient Φ , consider $\rho \in L^\gamma(\mathbb{T}^N)$, $\rho \geq 0$, and $\mathbf{v} \in L^2(\mathbb{T}^N)$ such that $\sqrt{\rho} \mathbf{v} \in L^2(\mathbb{T}^N)$. Denote $\mathbf{q} = \rho \mathbf{v}$ and let $\Phi(\rho, \mathbf{q}) : \mathfrak{U} \rightarrow L^1(\mathbb{T}^N)$ be defined as follows

$$\Phi(\rho, \mathbf{q}) e_k = \mathbf{g}_k(\cdot, \rho(\cdot), \mathbf{q}(\cdot)) = \mathbf{h}_k(\cdot, \rho(\cdot)) + \alpha_k \mathbf{q}(\cdot),$$

where the coefficients $\alpha_k \in \mathbb{R}$ are constants and $\mathbf{h}_k : \mathbb{T}^N \times \mathbb{R} \rightarrow \mathbb{R}$ are C^1 -functions that satisfy

$$(2.2) \quad \sum_{k \geq 1} |\alpha_k|^2 < \infty,$$

$$(2.3) \quad \sum_{k \geq 1} |\mathbf{h}_k(x, \rho)|^2 \leq C(\rho^2 + |\rho|^{\gamma+1}),$$

$$(2.4) \quad \sum_{k \geq 1} |\nabla_\rho \mathbf{h}_k(x, \rho)|^2 \leq C(1 + |\rho|^{\gamma-1}).$$

Remark that in this setting $L^1(\mathbb{T}^N)$ is the natural space for values of the operator $\Phi(\rho, \rho \mathbf{v})$. Indeed, due to lack of a priori estimates for (1.1) it is not possible to consider $\Phi(\rho, \rho \mathbf{v})$ as a mapping with values in a space with higher integrability. This fact brings difficulties concerning the definition of the stochastic integral in (1.1) because the space $L^1(\mathbb{T}^N)$ does not belong among 2-smooth Banach spaces nor among UMD Banach spaces where the theory of stochastic Itô integration is well-established (see e.g. [3], [21], [19]). However, since we expect the momentum equation (1.1b) to be satisfied only in the sense of distributions anyway, we make use of the embedding $L^1(\mathbb{T}^N) \hookrightarrow W^{-l,2}(\mathbb{T}^N)$, which is true provided $l > \frac{N}{2}$, and understand the stochastic integral as a process in the Hilbert space $W^{-l,2}(\mathbb{T}^N)$. To be more precise, it is easy to check that under the above assumptions on ρ and \mathbf{v} , the mapping $\Phi(\rho, \rho \mathbf{v})$ belongs to $L_2(\mathfrak{U}; W^{-l,2}(\mathbb{T}^N))$, the space of Hilbert-Schmidt operators from \mathfrak{U} to $W^{-l,2}(\mathbb{T}^N)$. Indeed, due to (2.2) and (2.3)

$$(2.5) \quad \begin{aligned} \|\Phi(\rho, \rho \mathbf{v})\|_{L_2(\mathfrak{U}; W_x^{-l,2})}^2 &= \sum_{k \geq 1} \|\mathbf{g}_k(\rho, \rho \mathbf{v})\|_{W_x^{-l,2}}^2 \leq C \sum_{k \geq 1} \|\mathbf{g}_k(\rho, \rho \mathbf{v})\|_{L_x^1}^2 \\ &\leq \sum_{k \geq 1} \left(\int_{\mathbb{T}^N} (|\mathbf{h}_k(x, \rho)| + \rho |\alpha_k \mathbf{v}|) dx \right)^2 \\ &\leq C(\rho)_{\mathbb{T}^N} \int_{\mathbb{T}^N} \left(\sum_{k \geq 1} \rho^{-1} |\mathbf{h}_k(x, \rho)|^2 + \sum_{k \geq 1} \rho |\alpha_k \mathbf{v}|^2 \right) dx \\ &\leq C(\rho)_{\mathbb{T}^N} \int_{\mathbb{T}^N} (\rho + \rho^\gamma + \rho |\mathbf{v}|^2) dx < \infty, \end{aligned}$$

where $(\rho)_{\mathbb{T}^N}$ denotes the mean value of ρ over \mathbb{T}^N . Consequently, if

$$\begin{aligned} \rho &\in L^\gamma(\Omega \times (0, T), \mathcal{P}, d\mathbb{P} \otimes dt; L^\gamma(\mathbb{T}^N)), \\ \sqrt{\rho} \mathbf{v} &\in L^2(\Omega \times (0, T), \mathcal{P}, d\mathbb{P} \otimes dt; L^2(\mathbb{T}^N)), \end{aligned}$$

where \mathcal{P} denotes the progressively measurable σ -algebra associated to (\mathcal{F}_t) , and the mean value $(\rho(t))_{\mathbb{T}^N}$ (that is constant in t but in general depends on ω) is for instance essentially bounded then the stochastic integral $\int_0^\cdot \Phi(\rho, \rho \mathbf{v}) dW$ is a well-defined (\mathcal{F}_t) -martingale taking values in $W^{-l,2}(\mathbb{T}^N)$. Finally, we define the auxiliary space $\mathfrak{U}_0 \supset \mathfrak{U}$ via

$$\mathfrak{U}_0 = \left\{ v = \sum_{k \geq 1} c_k e_k; \sum_{k \geq 1} \frac{c_k^2}{k^2} < \infty \right\},$$

endowed with the norm

$$\|v\|_{\mathfrak{U}_0}^2 = \sum_{k \geq 1} \frac{c_k^2}{k^2}, \quad v = \sum_{k \geq 1} c_k e_k.$$

Note that the embedding $\mathfrak{U} \hookrightarrow \mathfrak{U}_0$ is Hilbert-Schmidt. Moreover, trajectories of W are \mathbb{P} -a.s. in $C([0, T]; \mathfrak{U}_0)$ (see [7]).

2.1. The concept of solution and the main result. Existence of the so-called finite energy weak martingale solution to the stochastic Navier-Stokes system for compressible fluids, in particular (1.1), was recently established in [2]. Let us recall the corresponding definition of a solution and the existence result.

Definition 2.1. Let Λ be a Borel probability measure on $L^\gamma(\mathbb{T}^N) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)$. Then

$$((\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}), \varrho, \mathbf{u}, W)$$

is called a finite energy weak martingale solution to (1.1) with the initial data Λ provided

- (a) $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ is a stochastic basis with a complete right-continuous filtration,
- (b) W is an (\mathcal{F}_t) -cylindrical Wiener process,
- (c) the density $\varrho \geq 0$ is (\mathcal{F}_t) -adapted and $\varrho \in L^\gamma(\Omega; C_w([0, T]; L^\gamma(\mathbb{T}^N)))$,
- (d) the velocity \mathbf{u} is (\mathcal{F}_t) -adapted and $\mathbf{u} \in L^2(\Omega; L^2(0, T; W^{1,2}(\mathbb{T}^N)))$,
- (e) the momentum $\varrho \mathbf{u} \in L^{\frac{2\gamma}{\gamma+1}}(\Omega; C_w([0, T]; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)))$,
- (f) $\Lambda = \mathbb{P} \circ (\varrho(0), \varrho \mathbf{u}(0))^{-1}$.
- (g) $\Phi(\varrho, \varrho \mathbf{u}) \in L^2(\Omega \times [0, T], \mathcal{P}, d\mathbb{P} \otimes dt; L_2(\mathfrak{U}; W^{-l,2}(\mathbb{T}^N)))$ for some $l > \frac{N}{2}$,
- (h) for all $\psi \in C^\infty(\mathbb{T}^N)$ and $\varphi \in C^\infty(\mathbb{T}^N)$ and all $t \in [0, T]$ it holds \mathbb{P} -a.s.

$$\begin{aligned} \langle \varrho(t), \psi \rangle &= \langle \varrho(0), \psi \rangle + \int_0^t \langle \varrho \mathbf{u}, \nabla \psi \rangle ds, \\ \langle \varrho \mathbf{u}(t), \varphi \rangle &= \langle \varrho \mathbf{u}(0), \varphi \rangle + \int_0^t \langle \varrho \mathbf{u} \otimes \mathbf{u}, \nabla \varphi \rangle ds - \nu \int_0^t \langle \nabla \mathbf{u}, \nabla \varphi \rangle ds \\ &\quad - (\lambda + \nu) \int_0^t \langle \operatorname{div} \mathbf{u}, \operatorname{div} \varphi \rangle ds + \frac{1}{\varepsilon^2} \int_0^t \langle \rho^\gamma, \operatorname{div} \varphi \rangle ds \\ &\quad + \int_0^t \langle \Phi(\varrho, \varrho \mathbf{u}) dW, \varphi \rangle, \end{aligned}$$

- (i) for all $p \in [1, \infty)$ the following energy inequality holds true

$$\begin{aligned} (2.6) \quad & \mathbb{E} \left[\sup_{0 \leq t \leq T} \int_{\mathbb{T}^N} \left(\frac{1}{2} \varrho(t) |\mathbf{u}(t)|^2 + \frac{1}{\varepsilon^2(\gamma-1)} \varrho^\gamma(t) \right) dx \right]^p \\ & + \mathbb{E} \left[\int_0^T \int_{\mathbb{T}^N} \nu |\nabla \mathbf{u}|^2 + (\lambda + \nu) |\operatorname{div} \mathbf{u}|^2 \right]^p \\ & \leq C(p) \mathbb{E} \left[\int_{\mathbb{T}^N} \left(\frac{1}{2} \frac{|\varrho \mathbf{u}(0)|^2}{\varrho(0)} + \frac{1}{\varepsilon^2(\gamma-1)} \varrho(0)^\gamma \right) dx + 1 \right]^p. \end{aligned}$$

- (j) Let $b \in C^1(\mathbb{R})$ such that $b'(z) = 0$ for all $z \geq M_b$. Then for all $\psi \in C^\infty(\mathbb{T}^N)$ and all $t \in [0, T]$ it holds \mathbb{P} -a.s.

$$\langle b(\varrho(t)), \psi \rangle = \langle b(\varrho(0)), \psi \rangle + \int_0^t \langle b(\varrho) \mathbf{u}, \nabla \psi \rangle ds - \int_0^t \langle (b'(\varrho) \varrho - b(\varrho) \mathbf{u}) \operatorname{div} \mathbf{u}, \psi \rangle ds.$$

Remark 2.2. In Def. 2.1 (j) the continuity equation is stated in the renormalized sense. This is part of the existence result in [2] but will not be used in the remainder of the paper.

Theorem 2.3. Assume that for the initial law Λ there exists $M \in (0, \infty)$ such that

$$\Lambda \left\{ (\rho, \mathbf{q}) \in L^\gamma(\mathbb{T}^N) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N); \rho \geq 0, (\rho)_{\mathbb{T}^N} \leq M, \mathbf{q}(x) = 0 \text{ if } \rho(x) = 0 \right\} = 1,$$

and that for all $p \in [1, \infty)$ the following moment estimate holds true

$$\int_{L_x^\gamma \times L_x^{\frac{2\gamma}{\gamma+1}}} \left\| \frac{1}{2} \frac{|\mathbf{q}|^2}{\rho} + \frac{1}{\varepsilon^2(\gamma-1)} \rho^\gamma \right\|_{L_x^1}^p d\Lambda(\rho, \mathbf{q}) \leq C_\varepsilon.$$

Then there exists a finite energy weak martingale solution to (1.1) with the initial data Λ .

Concerning the incompressible Navier-Stokes system (1.2), several notions of solution are typically considered depending on the space dimension. From the PDE point of view, we restrict ourselves to weak solutions (although more can be proved in dimension two), i.e. (1.2) is satisfied in the sense of distributions. From the probabilistic point of view, we will consider two concepts, namely, pathwise (or strong) solutions and martingale (or weak) solutions. In the former one the underlying probability space as well as the driving process is fixed in advance while in the latter case these stochastic elements become part of the solution of the problem. Clearly, existence of a pathwise solution is stronger and implies existence of a martingale solution. Besides, due to classical Yamada-Watanabe-type argument (see e.g. [13], [22]), existence of a pathwise solution follows from existence of a martingale solution together with pathwise uniqueness. The difference lies also in the way how the initial condition is posed: for pathwise solutions we are given a random variable \mathbf{u}_0 whereas for martingale solutions we can only prescribe an initial law Λ .

Note that due to our assumptions on the operator Φ , the stochastic perturbations that we obtain in the limit system (1.2) is affine linear function of the velocity and takes the following form

$$\Psi(\mathbf{v})e_k d\beta_k = \mathcal{P}_H \Phi(1, \mathbf{v})e_k d\beta_k = (\mathcal{P}_H \mathbf{h}_k(1) + \alpha_k \mathbf{v})d\beta_k.$$

Besides, due to (2.2), (2.3) it holds true that

$$(2.7) \quad \begin{aligned} \|\Psi(\mathbf{v})\|_{L^2(\mathfrak{U}; L_x^2)}^2 &\leq C(1 + \|\mathbf{v}\|_{L_x^2}^2), \\ \|\Psi(\mathbf{v}) - \Psi(\mathbf{w})\|_{L^2(\mathfrak{U}; L_x^2)}^2 &\leq C\|\mathbf{v} - \mathbf{w}\|_{L_x^2}^2. \end{aligned}$$

In dimension three, existence of a strong solution which is closely related to uniqueness is one the celebrated Millenium Prize Problems and remains unsolved. Therefore, we consider weak martingale solutions, see for instance [6] or [12].

Definition 2.4. *Let Λ be a Borel probability measure on $L^2(\mathbb{T}^N)$. Then*

$$((\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P}), \mathbf{u}, W)$$

is called a weak martingale solution to (1.2) with the initial data Λ provided

- (a) $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ is a stochastic basis with a complete right-continuous filtration,
- (b) W is an (\mathcal{F}_t) -cylindrical Wiener process,
- (c) the velocity \mathbf{u} is (\mathcal{F}_t) -adapted and

$$\mathbf{u} \in L^2(\Omega; L^2(0, T; W_{div}^{1,2}(\mathbb{T}^N))) \cap L^2(\Omega; C_w([0, T]; L_{div}^2(\mathbb{T}^N))),$$

- (d) $\Lambda = \mathbb{P} \circ \mathbf{u}(0)^{-1}$,
- (e) for all $\varphi \in C_{div}^\infty(\mathbb{T}^N)$ and all $t \in [0, T]$ it holds \mathbb{P} -a.s.

$$\langle \mathbf{u}(t), \varphi \rangle = \langle \mathbf{u}(0), \varphi \rangle + \int_0^t \langle \mathbf{u} \otimes \mathbf{u}, \nabla \varphi \rangle ds - \nu \int_0^t \langle \nabla \mathbf{u}, \nabla \varphi \rangle ds + \int_0^t \langle \Psi(\mathbf{u}) dW, \varphi \rangle.$$

Here and hereafter, the substrict div refers to the space of solenoidal (divergenceless) functions.

Under the condition (2.7), the following existence result holds true and can be found for instance in [6] and [12].

Theorem 2.5. *Let Λ be a Borel probability measure on $L^2(\mathbb{T}^N)$ such that for all $p \in [1, \infty)$*

$$\int_{L_x^2} \|\mathbf{v}\|_{L_x^2}^p d\Lambda(\mathbf{v}) \leq C(p).$$

Then there exists a weak martingale solution to (1.2) with initial law Γ .

In dimension two, pathwise uniqueness for weak solutions is known under (2.7), we refer the reader for instance to [5], [4]. Consequently, we may work with the definition of a weak pathwise solution.

Definition 2.6. Let $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ be a given stochastic basis with an (\mathcal{F}_t) -cylindrical Wiener process W and let \mathbf{u}_0 be an \mathcal{F}_0 -measurable random variable. Then \mathbf{u} is called a weak pathwise solution to (1.2) with the initial condition \mathbf{u}_0 provided

(a) the velocity \mathbf{u} is (\mathcal{F}_t) -adapted and

$$\mathbf{u} \in L^2(\Omega; L^2(0, T; W_{div}^{1,2}(\mathbb{T}^N))) \cap L^2(\Omega; C_w([0, T]; L_{div}^2(\mathbb{T}^N))),$$

(b) $\mathbf{u}(0) = \mathbf{u}_0$ \mathbb{P} -a.s.,

(c) for all $\varphi \in C_{div}^\infty(\mathbb{T}^N)$ and all $t \in [0, T]$ it holds \mathbb{P} -a.s.

$$\langle \mathbf{u}(t), \varphi \rangle = \langle \mathbf{u}_0, \varphi \rangle + \int_0^t \langle \mathbf{u} \otimes \mathbf{u}, \nabla \varphi \rangle ds - \nu \int_0^t \langle \nabla \mathbf{u}, \nabla \varphi \rangle ds + \int_0^t \langle \Psi(\mathbf{u}) dW, \varphi \rangle.$$

Theorem 2.7. Let $N = 2$. Let $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ be a given stochastic basis with an (\mathcal{F}_t) -cylindrical Wiener process W and let \mathbf{u}_0 be an \mathcal{F}_0 -measurable random variable such that $\mathbf{u}_0 \in L^p(\Omega; L^2(\mathbb{T}^2))$ for all $p \in [1, \infty)$. Then there exists a unique weak pathwise solution to (1.2) with the initial condition \mathbf{u}_0 .

The main results of the present paper are following.

Theorem 2.8. Let Λ be a given Borel probability measure on $L^2(\mathbb{T}^N)$. Let Λ_ε be a Borel probability measure on $L^\gamma(\mathbb{T}^N) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)$ such that for some constant $M > 0$ (independent of ε) it holds true that

$$\Lambda_\varepsilon \left\{ (\rho, \mathbf{q}) \in L^\gamma(\mathbb{T}^N) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N); \rho \geq \frac{1}{M}, \left| \frac{\rho - 1}{\varepsilon} \right| \leq M \right\} = 1,$$

for all $p \in [1, \infty)$,

$$\int_{L_x^\gamma \times L_x^{\frac{2\gamma}{\gamma+1}}} \left\| \frac{1}{2} \frac{|\mathbf{q}|^2}{\rho} \right\|_{L_x^p}^p d\Lambda_\varepsilon(\rho, \mathbf{q}) \leq C(p),$$

and that the marginal law of Λ_ε corresponding to the second component converges to Λ weakly in the sense of measures on $L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)$. If $((\Omega^\varepsilon, \mathcal{F}^\varepsilon, (\mathcal{F}^\varepsilon), \mathbb{P}^\varepsilon), \varrho_\varepsilon, \mathbf{u}_\varepsilon, W_\varepsilon)$ is a finite energy weak martingale solution to (1.1) with the initial law Λ_ε , $\varepsilon \in (0, 1)$, then¹

$$\begin{aligned} \varrho_\varepsilon &\rightarrow 1 \quad \text{in law on } L^\infty(0, T; L^\gamma(\mathbb{T}^N)), \\ \mathbf{u}_\varepsilon &\rightarrow \mathbf{u} \quad \text{in law on } (L^2(0, T; W^{1,2}(\mathbb{T}^N)), w), \end{aligned}$$

where \mathbf{u} is a weak martingale solution to (1.2) with the initial law Λ .

Theorem 2.9. Let $N = 2$ and $\mathbf{u}_0 \in L^2(\mathbb{T}^2)$. Let Λ_ε be a Borel probability measure on $L^\gamma(\mathbb{T}^2) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^2)$ such that for some constant $M > 0$ (independent of ε) it holds true that

$$\Lambda_\varepsilon \left\{ (\rho, \mathbf{q}) \in L^\gamma(\mathbb{T}^2) \times L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^2); \rho \geq \frac{1}{M}, \left| \frac{\rho - 1}{\varepsilon} \right| \leq M, \left| \frac{\mathbf{q} - \mathbf{u}_0}{\varepsilon} \right| \leq M \right\} = 1,$$

If $((\Omega, \mathcal{F}, (\mathcal{F}), \mathbb{P}), \varrho_\varepsilon, \mathbf{u}_\varepsilon, W)$ is a finite energy weak martingale solution to (1.1) with the initial law Λ_ε , $\varepsilon \in (0, 1)$, then

$$\begin{aligned} \varrho_\varepsilon &\rightarrow 1 \quad \text{in } L^\infty(0, T; L^\gamma(\mathbb{T}^2)) \quad \mathbb{P}\text{-a.s.}, \\ \mathbf{u}_\varepsilon &\rightarrow \mathbf{u} \quad \text{in } (L^2(0, T; W^{1,2}(\mathbb{T}^2)), w) \quad \mathbb{P}\text{-a.s.}, \end{aligned}$$

where \mathbf{u} is a weak pathwise solution to (1.2) with the initial condition \mathbf{u}_0 .

¹If a topological space X is equipped with the weak topology we write (X, w) .

Here and in the sequel, the letter C denotes a constant that might change from one line to another and that is independent of ε .

3. PROOF OF THEOREM 2.8

This section is devoted to the study the limit $\varepsilon \rightarrow 0$ in the system (1.1). To this end, we recall that it was proved in [2] that for every $\varepsilon \in (0, 1)$ there exists

$$((\Omega^\varepsilon, \mathcal{F}^\varepsilon, (\mathcal{F}_t^\varepsilon), \mathbb{P}^\varepsilon), \varrho_\varepsilon, \mathbf{u}_\varepsilon, W_\varepsilon)$$

which is a weak martingale solution in the sense of Definition 2.1. It was shown in [14] that it is enough to consider only one probability space, namely,

$$(\Omega^\varepsilon, \mathcal{F}^\varepsilon, \mathbb{P}^\varepsilon) = ([0, 1], \mathcal{B}([0, 1]), \mathcal{L}) \quad \forall \varepsilon \in (0, 1)$$

where \mathcal{L} denotes the Lebesgue measure on $[0, 1]$. Moreover, we can assume without loss of generality that there exists one common Wiener process W for all ε .

3.1. Uniform bounds. We start with an a priori estimate which is a modification of the energy estimate (2.6) established in [2].

Proposition 3.1. *Let $p \in [1, \infty)$. Then the following estimate holds true uniformly in ε*

$$\begin{aligned} & \mathbb{E} \left[\sup_{0 \leq t \leq T} \int_{\mathbb{T}^N} \left(\frac{1}{2} \varrho_\varepsilon(t) |\mathbf{u}_\varepsilon(t)|^2 + \frac{1}{\varepsilon^2(\gamma - 1)} (\varrho_\varepsilon^\gamma(t) - 1 - \gamma(\varrho_\varepsilon(t) - 1)) \right) dx \right]^p \\ (3.1) \quad & + \mathbb{E} \left[\int_0^T \int_{\mathbb{T}^N} \nu |\nabla \mathbf{u}_\varepsilon|^2 + (\lambda + \nu) |\operatorname{div} \mathbf{u}_\varepsilon|^2 \right]^p \\ & \leq C_p \mathbb{E} \left[\int_{\mathbb{T}^N} \left(\frac{1}{2} \varrho_\varepsilon(0) |\mathbf{u}_\varepsilon(0)|^2 + \frac{1}{\varepsilon^2(\gamma - 1)} (\varrho_\varepsilon^\gamma(0) - 1 - \gamma(\varrho_\varepsilon(0) - 1)) \right) dx + 1 \right]^p \\ & \leq C_p. \end{aligned}$$

Proof. The first inequality follows directly from Definition 2.1 and the mass conservation

$$\int_{\mathbb{T}^N} \varrho_\varepsilon(t) dx = \int_{\mathbb{T}^N} \varrho_\varepsilon(0) dx$$

which is a consequence of equation (1.1a). Next, we observe that due to the Taylor theorem and our assumptions upon Λ_ε , it holds

$$\mathbb{E} \int_{\mathbb{T}^N} (\varrho_\varepsilon^\gamma(0) - 1 - \gamma(\varrho_\varepsilon(0) - 1)) dx \leq C\varepsilon^2$$

and hence the second estimate follows (independently of ε). \square

Consequently, we gain the uniform bounds, for all $p \in [1, \infty)$,

$$(3.2) \quad \nabla \mathbf{u}_\varepsilon \in L^p(\Omega; L^2(0, T; L^2(\mathbb{T}^N))),$$

$$(3.3) \quad \sqrt{\varrho_\varepsilon} \mathbf{u}_\varepsilon \in L^p(\Omega; L^\infty(0, T; L^2(\mathbb{T}^N))).$$

Moreover, as

$$\int_{\mathbb{T}^N} \varrho_\varepsilon(t) dx = \int_{\mathbb{T}^N} \varrho_\varepsilon(0) dx \geq \frac{1}{M} |\mathbb{T}^N|,$$

the above estimates give rise to

$$(3.4) \quad \mathbf{u}_\varepsilon \in L^p(\Omega; L^2(0, T; W^{1,2}(\mathbb{T}^N))).$$

Let us now introduce the essential and residual component of any function h :

$$h = h_{\text{ess}} + h_{\text{res}},$$

$$h_{\text{ess}} = \chi(\varrho_\varepsilon)h, \quad \chi \in C_c^\infty(0, \infty), \quad 0 \leq \chi \leq 1, \quad \chi = 1 \text{ on an open interval containing } 1,$$

$$h_{\text{res}} = (1 - \chi(\varrho_\varepsilon))h.$$

The following lemma will be useful.

Lemma 3.2. *Let $P(\rho) := \rho^\gamma - 1 - \gamma(\rho - 1)$, with $\rho \in [0, \infty)$. Then there exist constants $C_1, C_2, C_3, C_4 > 0$ such that*

- (i) $C_1|\rho - 1|^2 \leq P(\rho) \leq C_2|\rho - 1|^2$ if $\rho \in \text{supp } \chi$,
- (ii) $P(\rho) \geq C_4$ if $\rho \notin \text{supp } \chi$,
- (iii) $P(\rho) \geq C_3\rho^\gamma$ if $\rho \notin \text{supp } \chi$.

Proof. The first statement follows immediately from the Taylor theorem. The second one is a consequence of the fact that P is strictly convex and attains its minimum at $\rho = 1$. If $\rho \notin \text{supp } \chi$ and $\rho \in [0, 1)$ then the third statement is a consequence of the second one. Finally, we observe that the function $\frac{P(\rho)}{\rho^\gamma}$ is increasing for large $\rho \in [1, \infty)$ and its value at $\rho = 1$ is zero. This implies the remaining part of (iii) and the proof is complete. \square

Accordingly, we obtain the following uniform bounds, for all $p \in [1, \infty)$

$$\begin{aligned} \left[\frac{\varrho_\varepsilon - 1}{\varepsilon} \right]_{\text{ess}} &\in L^p(\Omega; L^\infty(0, T; L^2(\mathbb{T}^N))), \\ \frac{[\varrho_\varepsilon]_{\text{res}} + [1]_{\text{res}}}{\varepsilon^2} &\in L^p(\Omega; L^\infty(0, T; L^\gamma(\mathbb{T}^N))), \end{aligned}$$

therefore, setting $\varphi_\varepsilon := \frac{1}{\varepsilon}(\varrho_\varepsilon - 1)$, we deduce that

$$(3.5) \quad \varphi_\varepsilon \in L^p(\Omega; L^\infty(0, T; L^{\min(\gamma, 2)}(\mathbb{T}^N))).$$

As the next step, we want to show that

$$(3.6) \quad \varrho_\varepsilon \rightarrow 1 \quad \text{in} \quad L^p(\Omega; L^\infty(0, T; L^\gamma(\mathbb{T}^N))),$$

which in particular leads to

$$(3.7) \quad \varrho_\varepsilon \in L^p(\Omega; L^\infty(0, T; L^\gamma(\mathbb{T}^N))).$$

Then, combining (3.3), (3.7) and (3.4) and (3.7), respectively, we deduce the uniform bounds, for all $p \in [1, \infty)$,

$$(3.8) \quad \varrho_\varepsilon \mathbf{u}_\varepsilon \in L^p(\Omega; L^\infty(0, T; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N))),$$

$$(3.9) \quad \varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon \in L^p(\Omega; L^2(0, T; L^{\frac{6\gamma}{4\gamma+3}}(\mathbb{T}^N))).$$

Let us now verify (3.6). Since for all $\delta > 0$ there exists $C_\delta > 0$ such that

$$\rho^\gamma - 1 - \gamma(\rho - 1) \geq C_\delta |\rho - 1|^\gamma$$

if $|\rho - 1| \geq \delta$ and $\rho \geq 0$, we obtain

$$\begin{aligned} \mathbb{E} \left[\sup_{0 \leq t \leq T} \int_{\mathbb{T}^N} |\varrho_\varepsilon - 1|^\gamma \, dx \, dt \right]^p &= \mathbb{E} \left[\sup_{0 \leq t \leq T} \int_{\mathbb{T}^N} \mathbf{1}_{\{|\varrho_\varepsilon - 1| \geq \delta\}} |\varrho_\varepsilon - 1|^\gamma \, dx \, dt \right]^p \\ &\quad + \mathbb{E} \left[\sup_{0 \leq t \leq T} \int_{\mathbb{T}^N} \mathbf{1}_{\{|\varrho_\varepsilon - 1| < \delta\}} |\varrho_\varepsilon - 1|^\gamma \, dx \, dt \right]^p \\ &\leq C_\delta \mathbb{E} \left[\int_{\mathbb{T}^N} (\varrho_\varepsilon^\gamma - 1 - \gamma(\varrho_\varepsilon - 1)) \, dx \, dt \right]^p + C\delta^{\gamma p} \leq C_\delta \varepsilon^{2p} + C\delta^{\gamma p}. \end{aligned}$$

Letting first $\varepsilon \rightarrow 0$ and then $\delta \rightarrow 0$ yields the claim.

3.2. Acoustic equation. In order to proceed we need the Helmholtz projection \mathcal{P}_H which projects $L^2(\mathbb{T}^N)$ onto divergence free vector fields

$$L_{\text{div}}^2(\mathbb{T}^N) := \overline{C_{\text{div}}^\infty(\mathbb{T}^N)}^{\|\cdot\|^2}.$$

Moreover, we set $\mathcal{Q} = \text{Id} - \mathcal{P}_H$. Recall that \mathcal{P}_H can be easily defined in terms of the Fourier coefficients a_k (cf. (2.1), in particular it can be shown that both \mathcal{P}_H and \mathcal{Q} are continuous in all $W^{l,q}(\mathbb{T}^N)$ -spaces, $l \in \mathbb{R}$, $q \in (1, \infty)$).

Let us now project (1.1b) onto the space of gradient vector fields. Then (1.1) rewrites as

$$(3.10a) \quad \varepsilon \, d\varphi_\varepsilon + \text{div} \, \mathcal{Q}(\varrho_\varepsilon \mathbf{u}_\varepsilon) dt = 0,$$

$$(3.10b) \quad \varepsilon \, d\mathcal{Q}(\varrho_\varepsilon \mathbf{u}_\varepsilon) + \gamma \nabla \varphi_\varepsilon dt = \varepsilon \mathbf{F}_\varepsilon dt + \varepsilon \mathcal{Q}\Phi(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) dW,$$

$$\mathbf{F}_\varepsilon = \nu \Delta \mathcal{Q} \mathbf{u}_\varepsilon + (\lambda + \nu) \nabla \text{div} \, \mathbf{u}_\varepsilon - \mathcal{Q}[\text{div}(\varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon)] - \frac{1}{\varepsilon^2} \nabla[\varrho_\varepsilon^\gamma - 1 - \gamma(\varrho_\varepsilon - 1)].$$

The system (3.10) may be viewed as a stochastic version of Lighthill's acoustic analogy associated to the compressible Navier-Stokes system. Note that Proposition 3.1 yields

$$(3.11) \quad \mathbf{F}_\varepsilon \in L^p(0, T; L^1(0, T; W^{-l,2}(\mathbb{T}^N)))$$

uniformly in ε .

3.3. Compactness. Let us define the path space $\mathcal{X} = \mathcal{X}_\varrho \times \mathcal{X}_\mathbf{u} \times \mathcal{X}_{\varrho\mathbf{u}} \times \mathcal{X}_W$ where

$$\begin{aligned} \mathcal{X}_\varrho &= C_w(0, T; L^\gamma(\mathbb{T}^N)), & \mathcal{X}_\mathbf{u} &= (L^2(0, T; W^{1,2}(\mathbb{T}^N)), w), \\ \mathcal{X}_{\varrho\mathbf{u}} &= C_w([0, T]; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)), & \mathcal{X}_W &= C([0, T]; \mathfrak{U}_0). \end{aligned}$$

Let us denote by $\mu_{\varrho_\varepsilon}$, $\mu_{\mathbf{u}_\varepsilon}$ and $\mu_{\mathcal{P}(\varrho_\varepsilon \mathbf{u}_\varepsilon)}$, respectively, the law of ϱ_ε , \mathbf{u}_ε , $\mathcal{P}(\varrho_\varepsilon \mathbf{u}_\varepsilon)$ on the corresponding path space. By μ_W we denote the law of W on \mathcal{X}_W and their joint law on \mathcal{X} is denoted by μ^ε .

To proceed, it is necessary to establish tightness of $\{\mu^\varepsilon; \varepsilon \in (0, 1)\}$.

Proposition 3.3. *The set $\{\mu_{\mathbf{u}_\varepsilon}; \varepsilon \in (0, 1)\}$ is tight on $\mathcal{X}_\mathbf{u}$.*

Proof. This is a consequence of (3.4). Indeed, for any $R > 0$ the set

$$B_R = \{\mathbf{u} \in L^2(0, T; W^{1,2}(\mathbb{T}^N)); \|\mathbf{u}\|_{L^2(0, T; W^{1,2}(\mathbb{T}^N))} \leq R\}$$

is relatively compact in $\mathcal{X}_\mathbf{u}$ and

$$\mu_{\mathbf{u}_\varepsilon}(B_R^c) = \mathbb{P}(\|\mathbf{u}_\varepsilon\|_{L^2(0, T; W^{1,2}(\mathbb{T}^N))} \geq R) \leq \frac{1}{R} \mathbb{E} \|\mathbf{u}_\varepsilon\|_{L^2(0, T; W^{1,2}(\mathbb{T}^N))} \leq \frac{C}{R}$$

which yields the claim. \square

Proposition 3.4. *The set $\{\mu_{\varrho_\varepsilon}; \varepsilon \in (0, 1)\}$ is tight on \mathcal{X}_ϱ .*

Proof. Due to (3.8), $\{\text{div}(\varrho_\varepsilon \mathbf{u}_\varepsilon)\}$ is bounded in $L^p(\Omega; L^\infty(0, T; W^{-1, \frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)))$ and therefore the continuity equation yields the following uniform bound, for all $p \in [1, \infty)$,

$$\varrho_\varepsilon \in L^p(\Omega; C^{0,1}([0, T]; W^{-1, \frac{2\gamma}{\gamma+1}}(\mathbb{T}^N))).$$

Now, the required tightness in follows by a similar reasoning as in Proposition 3.3 together with (3.7) and the compact embedding (see [20, Corollary B.2])

$$L^\infty(0, T; L^\gamma(\mathbb{T}^N)) \cap C^{0,1}([0, T]; W^{-2, \frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)) \xrightarrow{c} C_w([0, T]; L^\gamma(\mathbb{T}^N)).$$

\square

Proposition 3.5. *The set $\{\mu_{\mathcal{P}_H(\varrho_\varepsilon \mathbf{u}_\varepsilon)}; \varepsilon \in (0, 1)\}$ is tight on $\mathcal{X}_{\varrho\mathbf{u}}$.*

Proof. We decompose $\mathcal{P}_H(\varrho_\varepsilon \mathbf{u}_\varepsilon)$ into two parts, namely, $\mathcal{P}_H(\varrho_\varepsilon \mathbf{u}_\varepsilon)(t) = Y^\varepsilon(t) + Z^\varepsilon(t)$, where

$$\begin{aligned} Y^\varepsilon(t) &= \mathcal{P}_H \mathbf{q}_\varepsilon(0) - \int_0^t \mathcal{P}_H [\operatorname{div}(\varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon) - \nu \Delta \mathbf{u}_\varepsilon] ds, \\ Z^\varepsilon(t) &= \int_0^t \mathcal{P}_H \Phi(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) dW(s). \end{aligned}$$

Hölder continuity of (Y^ε) . We show that there exists $l \in \mathbb{N}$ such that for all $\kappa \in (0, 1/2)$ it holds true

$$(3.12) \quad \mathbb{E} \|Y^\varepsilon\|_{C^\kappa([0, T]; W^{-l, 2}(\mathbb{T}^N))} \leq C.$$

Choose l such that $L^1(\mathbb{T}^N) \hookrightarrow W^{-l+1, 2}(\mathbb{T}^N)$. The a priori estimates (3.4) and (3.9) and the continuity of \mathcal{P} yield

$$\begin{aligned} \mathbb{E} \left\| Y^\varepsilon(t) - Y^\varepsilon(s) \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta &= \mathbb{E} \left\| \int_s^t \mathcal{P} [\operatorname{div}(\varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon) + \nu \Delta \mathbf{u}_\varepsilon] ds \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta \\ &\leq C \mathbb{E} \left\| \int_s^t \operatorname{div}(\varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon) ds \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta + C \mathbb{E} \left\| \int_s^t \Delta \mathbf{u}_\varepsilon ds \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta \\ &\leq C \mathbb{E} \left\| \int_s^t \varrho_\varepsilon \mathbf{u}_\varepsilon \otimes \mathbf{u}_\varepsilon ds \right\|_{L^1(\mathbb{T}^N)}^\theta + C \mathbb{E} \left\| \int_s^t \nabla \mathbf{u}_\varepsilon ds \right\|_{L^1(\mathbb{T}^N)}^\theta \leq C |t - s|^{\theta/2} \end{aligned}$$

and (3.12) follows by the Kolmogorov continuity criterion.

Hölder continuity of (Z^ε) . Next, we show that also

$$\mathbb{E} \|Z^\varepsilon\|_{C^\kappa([0, T]; W^{-l, 2}(\mathbb{T}^N))} \leq C,$$

where $l \in \mathbb{N}$ was given by the previous step and $\kappa \in (0, 1/2)$. From the embedding $L^1(\mathbb{T}^N) \hookrightarrow W^{-l, 2}(\mathbb{T}^N)$, (2.2), (2.3), the a priori estimates and the continuity of \mathcal{P}_H we get

$$\begin{aligned} \mathbb{E} \left\| Z^\varepsilon(t) - Z^\varepsilon(s) \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta &= \mathbb{E} \left\| \int_s^t \mathcal{P}_H \Phi(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) dW \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta \\ &\leq C \mathbb{E} \left\| \int_s^t \Phi(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) dW \right\|_{W^{-l, 2}(\mathbb{T}^N)}^\theta \leq C \mathbb{E} \left(\int_s^t \sum_{k \geq 1} \|\mathbf{g}^k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)\|_{W^{-l, 2}}^2 dr \right)^{\frac{\theta}{2}} \\ &\leq C \mathbb{E} \left(\int_s^t \sum_{k \geq 1} \|\mathbf{g}^k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)\|_{L^1}^2 dr \right)^{\frac{\theta}{2}} \\ &\leq C \mathbb{E} \left(\int_s^t \int_{\mathbb{T}^N} (\varrho_\varepsilon + \varrho_\varepsilon |\mathbf{u}_\varepsilon|^2 + \varrho_\varepsilon^\gamma) dx dr \right)^{\frac{\theta}{2}} \\ &\leq C |t - s|^{\frac{\theta}{2}} \left(1 + \mathbb{E} \sup_{0 \leq t \leq T} \|\sqrt{\varrho_\varepsilon} \mathbf{u}_\varepsilon\|_{L^2}^\theta + \mathbb{E} \sup_{0 \leq t \leq T} \|\varrho_\varepsilon\|_{L^\gamma}^{\theta\gamma/2} \right) \leq C |t - s|^{\frac{\theta}{2}} \end{aligned}$$

and the Kolmogorov continuity criterion applies.

Conclusion. Collecting the above results we obtain that

$$\mathbb{E} \|\mathcal{P}_H(\varrho_\varepsilon \mathbf{u}_\varepsilon)\|_{C^\kappa([0, T]; W^{-l, 2}(\mathbb{T}^N))} \leq C$$

for some $l \in \mathbb{N}$ and all $\kappa \in (0, 1/2)$. This implies the desired tightness by making use of (3.8), continuity of \mathcal{P}_H together with the compact embedding (see [20, Corollary B.2])

$$L^\infty(0, T; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)) \cap C^\kappa([0, T]; W^{-l, 2}(\mathbb{T}^N)) \xrightarrow{c} C_w([0, T]; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)).$$

□

Since also the law μ_W is tight as being Radon measures on the Polish space \mathcal{X}_W we can finally deduce tightness of the joint laws μ^ε .

Corollary 3.6. *The set $\{\mu^\varepsilon; \varepsilon \in (0, 1)\}$ is tight on \mathcal{X} .*

The path space \mathcal{X} is not a Polish space and so our compactness argument is based on the Jakubowski-Skorokhod representation theorem instead of the classical Skorokhod representation theorem, see [14]. To be more precise, passing to a weakly convergent subsequence μ^ε (and denoting by μ the limit law) we infer the following result.

Proposition 3.7. *There exists a subsequence μ^ε , a probability space $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$ with \mathcal{X} -valued Borel measurable random variables $(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\mathbf{q}}_\varepsilon, \tilde{W}_\varepsilon)$, $n \in \mathbb{N}$, and $(\tilde{\varrho}, \tilde{\mathbf{u}}, \tilde{\mathbf{q}}, \tilde{W})$ such that*

- (a) *the law of $(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\mathbf{q}}_\varepsilon, \tilde{W}_\varepsilon)$ is given by μ^ε , $\varepsilon \in (0, 1)$,*
- (b) *the law of $(\tilde{\varrho}, \tilde{\mathbf{u}}, \tilde{\mathbf{q}}, \tilde{W})$, denoted by μ , is a Radon measure,*
- (c) *$(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\mathbf{q}}_\varepsilon, \tilde{W}_\varepsilon)$ converges $\tilde{\mathbb{P}}$ -a.s. to $(\tilde{\varrho}, \tilde{\mathbf{u}}, \tilde{\mathbf{q}}, \tilde{W})$ in the topology of \mathcal{X} .*

Let us now fix some notation that will be used in the sequel. We denote by \mathbf{r}_t the operator of restriction to the interval $[0, t]$ acting on various path spaces. In particular, if X stands for one of the path spaces \mathcal{X}_ϱ , $\mathcal{X}_\mathbf{u}$ or \mathcal{X}_W and $t \in [0, T]$, we define

$$(3.13) \quad \mathbf{r}_t : X \rightarrow X|_{[0, t]}, \quad f \mapsto f|_{[0, t]}.$$

Clearly, \mathbf{r}_t is a continuous mapping. Let $(\tilde{\mathcal{F}}_t^\varepsilon)$ and $(\tilde{\mathcal{F}}_t)$, respectively, be the $\tilde{\mathbb{P}}$ -augmented canonical filtration of the process $(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{W}_\varepsilon)$ and $(\tilde{\varrho}, \tilde{\mathbf{u}}, \tilde{W})$, respectively, that is

$$\begin{aligned} \tilde{\mathcal{F}}_t^\varepsilon &= \sigma(\sigma(\mathbf{r}_t \tilde{\varrho}_\varepsilon, \mathbf{r}_t \tilde{\mathbf{u}}_\varepsilon, \mathbf{r}_t \tilde{W}_\varepsilon) \cup \{N \in \tilde{\mathcal{F}}; \tilde{\mathbb{P}}(N) = 0\}), \quad t \in [0, T], \\ \tilde{\mathcal{F}}_t &= \sigma(\sigma(\mathbf{r}_t \tilde{\varrho}, \mathbf{r}_t \tilde{\mathbf{u}}, \mathbf{r}_t \tilde{W}) \cup \{N \in \tilde{\mathcal{F}}; \tilde{\mathbb{P}}(N) = 0\}), \quad t \in [0, T]. \end{aligned}$$

3.4. Identification of the limit. The aim of this subsection is to identify the limit processes given by Proposition 3.7 with a weak martingale solution to (1.2). Namely, we prove the following result which in turn verifies Theorem 2.8.

Theorem 3.8. *The process \tilde{W} is a $(\tilde{\mathcal{F}}_t)$ -cylindrical Wiener process and*

$$((\tilde{\Omega}, \tilde{\mathcal{F}}, (\tilde{\mathcal{F}}_t), \tilde{\mathbb{P}}), \tilde{\mathbf{u}}, \tilde{W})$$

is a weak martingale solution to (1.2) with the initial law Λ .

The proof proceeds in several steps. First of all, we show that also on the new probability space $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$, the approximations $\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon$ solve the corresponding compressible Navier-Stokes system (1.1).

Proposition 3.9. *Let $\varepsilon \in (0, 1)$. The process \tilde{W}_ε is a $(\tilde{\mathcal{F}}_t)$ -cylindrical Wiener process and*

$$((\tilde{\Omega}, \tilde{\mathcal{F}}, (\tilde{\mathcal{F}}_t^\varepsilon), \tilde{\mathbb{P}}), \tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{W}_\varepsilon)$$

is a finite energy weak martingale solution to (1.1) with initial law Λ_ε .

Proof. The first part of the claim follows immediately from the fact that \tilde{W}_ε has the same law as W . As a consequence, there exists a collection of mutually independent real-valued $(\tilde{\mathcal{F}}_t)$ -Wiener processes $(\tilde{\beta}_k^\varepsilon)_{k \geq 1}$ such that $\tilde{W}_\varepsilon = \sum_{k \geq 1} \tilde{\beta}_k^\varepsilon e_k$.

To show that the continuity equation (1.1a) is satisfied, let us define, for all $t \in [0, T]$ and $\psi \in C^\infty(\mathbb{T}^N)$, the functional

$$L(\rho, \mathbf{q})_t = \langle \rho(t), \psi \rangle - \langle \rho(0), \psi \rangle - \int_0^t \langle \mathbf{q}, \nabla \psi \rangle ds.$$

Note that $(\rho, \mathbf{q}) \mapsto L(\rho, \mathbf{q})_t$ is continuous on $\mathcal{X}_\rho \times \mathcal{X}_{\mathbf{q}}$. Hence the laws of $L(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_t$ and $L(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_t$ coincide and since $(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)$ solves (1.1a) we deduce that

$$\tilde{\mathbb{E}}|L(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_t|^2 = \mathbb{E}|L(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_t|^2 = 0$$

hence $(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)$ solves (1.1a).

To verify the momentum equation (1.1b), we define for all $t \in [0, T]$ and $\varphi \in C^\infty(\mathbb{T}^N)$ the functionals

$$\begin{aligned} M(\rho, \mathbf{v}, \mathbf{q})_t &= \langle \mathbf{q}(t), \varphi \rangle - \langle \mathbf{q}(0), \varphi \rangle + \int_0^t \langle \mathbf{q} \otimes \mathbf{v}, \nabla \varphi \rangle dr - \nu \int_0^t \langle \nabla \mathbf{v}, \nabla \varphi \rangle dr \\ &\quad - (\lambda + \nu) \int_0^t \langle \operatorname{div} \mathbf{v}, \operatorname{div} \varphi \rangle dr + \frac{a}{\varepsilon^2} \int_0^t \langle \rho^\gamma, \operatorname{div} \varphi \rangle dr \\ N(\rho, \mathbf{q})_t &= \sum_{k \geq 1} \int_0^t \langle \mathbf{g}_k(\rho, \mathbf{q}), \varphi \rangle^2 dr, \\ N_k(\rho, \mathbf{q})_t &= \int_0^t \langle \mathbf{g}_k(\rho, \mathbf{q}), \varphi \rangle dr, \end{aligned}$$

let $M(\rho, \mathbf{v}, \mathbf{q})_{s,t}$ denote the increment $M(\rho, \mathbf{v}, \mathbf{q})_t - M(\rho, \mathbf{v}, \mathbf{q})_s$ and similarly for $N(\rho, \mathbf{q})_{s,t}$ and $N_k(\rho, \mathbf{q})_{s,t}$. We claim that with the above uniform estimates in hand, the mappings

$$(\rho, \mathbf{v}, \mathbf{q}) \mapsto M(\rho, \mathbf{v}, \mathbf{q})_t, \quad (\rho, \mathbf{v}, \mathbf{q}) \mapsto N(\rho, \mathbf{q})_t, \quad (\rho, \mathbf{v}, \mathbf{q}) \mapsto N_k(\rho, \mathbf{q})_t$$

are well-defined and measurable on a subspace of $\mathcal{X}_\rho \times \mathcal{X}_{\mathbf{v}} \times \mathcal{X}_{\mathbf{q}}$ where the joint law of $(\tilde{\varrho}, \tilde{\mathbf{u}}, \tilde{\varrho} \tilde{\mathbf{u}})$ is supported, i.e. where all the uniform estimates hold true. Indeed, in the case of $N(\rho, \mathbf{q})_t$ we have by (2.2), (2.3) similarly to (2.5)

$$\sum_{k \geq 1} \int_0^t \langle \mathbf{g}_k(\rho, \mathbf{q}), \varphi \rangle^2 ds \leq C \sum_{k \geq 1} \int_0^t \|\mathbf{g}_k(\rho, \mathbf{q})\|_{L^1}^2 ds \leq C.$$

$M(\rho, \mathbf{v}, \mathbf{q})$ and $N_k(\rho, \mathbf{v})_t$ can be handled similarly and therefore, the following random variables have the same laws

$$\begin{aligned} M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) &\stackrel{d}{\sim} M(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon), \\ N(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) &\stackrel{d}{\sim} N(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon), \\ N_k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) &\stackrel{d}{\sim} N_k(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon). \end{aligned}$$

Let us now fix times $s, t \in [0, T]$ such that $s < t$ and let

$$h : \mathcal{X}_\rho|_{[0,s]} \times \mathcal{X}_{\mathbf{u}}|_{[0,s]} \times \mathcal{X}_W|_{[0,s]} \rightarrow [0, 1]$$

be a continuous function. Since

$$M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_t = \int_0^t \langle \Phi(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) dW, \varphi \rangle = \sum_{k \geq 1} \int_0^t \langle \mathbf{g}_k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon), \varphi \rangle d\beta_k$$

is a square integrable (\mathcal{F}_t) -martingale, we infer that

$$[M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)]^2 - N(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon), \quad M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)\beta_k - N_k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)$$

are (\mathcal{F}_t) -martingales. Besides, it follows from the equality of laws that

$$\begin{aligned} (3.14) \quad &\tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\varrho}_\varepsilon, \mathbf{r}_s \tilde{\mathbf{u}}_\varepsilon, \mathbf{r}_s \tilde{W}_\varepsilon) [M(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_{s,t}] \\ &= \mathbb{E} h(\mathbf{r}_s \varrho_\varepsilon, \mathbf{r}_s \mathbf{u}_\varepsilon, \mathbf{r}_s W_\varepsilon) [M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_{s,t}] = 0, \end{aligned}$$

$$(3.15) \quad \begin{aligned} & \tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\varrho}_\varepsilon, \mathbf{r}_s \tilde{\mathbf{u}}_\varepsilon, \mathbf{r}_s \tilde{W}_\varepsilon) \left[[M(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^2]_{s,t} - N(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_{s,t} \right] \\ &= \mathbb{E} h(\mathbf{r}_s \varrho_\varepsilon, \mathbf{r}_s \mathbf{u}_\varepsilon, \mathbf{r}_s W_\varepsilon) \left[[M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)^2]_{s,t} - N(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_{s,t} \right] = 0, \end{aligned}$$

$$(3.16) \quad \begin{aligned} & \tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\varrho}_\varepsilon, \mathbf{r}_s \tilde{\mathbf{u}}_\varepsilon, \mathbf{r}_s \tilde{W}_\varepsilon) \left[[M(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \tilde{\beta}_k^\varepsilon]_{s,t} - N_k(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_{s,t} \right] \\ &= \mathbb{E} h(\mathbf{r}_s \varrho_\varepsilon, \mathbf{r}_s \mathbf{u}_\varepsilon, \mathbf{r}_s W_\varepsilon) \left[[M(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon) \beta_k]_{s,t} - N_k(\varrho_\varepsilon, \varrho_\varepsilon \mathbf{u}_\varepsilon)_{s,t} \right] = 0. \end{aligned}$$

The proof is hereby complete. \square

Consequently, we recover the result of Proposition 3.1 together with all the uniform estimates of the previous subsection. In particular, we find (for a subsequence) that

$$(3.17) \quad \tilde{\varrho}_\varepsilon \rightarrow 1 \quad \text{in } L^\infty(0, T; L^\gamma(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

Corollary 3.10. *We have the following bounds uniform in ε , for all $p \in [1, \infty)$ and $l > \frac{N}{2}$,*

$$\begin{aligned} \sqrt{\tilde{\varphi}_\varepsilon} \tilde{\mathbf{u}}_\varepsilon &\in L^p(\Omega, L^\infty(0, T; L^2(\mathbb{T}^N))), \\ \tilde{\varphi}_\varepsilon &\in L^p(\Omega, L^\infty(0, T; L^{\min(2, \gamma)}(\mathbb{T}^N))), \\ \tilde{\mathbf{F}}_\varepsilon &\in L^p(0, T; L^1(0, T; W^{-l, 2}(\mathbb{T}^N))) \end{aligned}$$

where $\tilde{\varphi}_\varepsilon = \frac{\tilde{\varrho}_\varepsilon - 1}{\varepsilon}$ and

$$\tilde{\mathbf{F}}_\varepsilon = \nu \Delta \mathbf{Q} \mathbf{u}_\varepsilon + (\lambda + \nu) \nabla \operatorname{div} \tilde{\mathbf{u}}_\varepsilon - \mathcal{Q}[\operatorname{div}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon \otimes \tilde{\mathbf{u}}_\varepsilon)] - \frac{1}{\varepsilon^2} \nabla[\tilde{\varrho}_\varepsilon^\gamma - 1 - \gamma(\tilde{\varrho}_\varepsilon - 1)].$$

Proposition 3.11. *We have the following convergence $\tilde{\mathbb{P}}$ -a.s.*

$$(3.18) \quad \mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon \rightarrow \tilde{\mathbf{u}} \quad \text{in } L^2(0, T; L^q(\mathbb{T}^N)) \quad \forall q < \frac{2N}{N-2}.$$

Proof. Since the joint laws of $(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \mathcal{P}_H(\varrho_\varepsilon \mathbf{u}_\varepsilon))$ and $(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\mathbf{q}}_\varepsilon)$ coincide, we deduce that $\tilde{\mathbf{q}}_\varepsilon = \mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)$ a.s. and consequently it follows from the proof of Proposition 3.5 that

$$(3.19) \quad \tilde{\mathbb{E}} \|\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)\|_{C^\kappa([0, T]; W^{-l, 2}(\mathbb{T}^N))} \leq C$$

for some $\kappa \in (0, 1)$ and $l \in \mathbb{N}$.

Besides, it follows from (3.17) and the convergence of $\tilde{\mathbf{u}}_\varepsilon$ to $\tilde{\mathbf{u}}$ that

$$(3.20) \quad \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon \rightharpoonup \tilde{\mathbf{u}} \quad \text{in } L^2(0, T; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

If we pass to the limit in the continuity equation, we see that $\operatorname{div} \tilde{\mathbf{u}} = 0$, which in turn identifies $\tilde{\mathbf{q}}$ with $\tilde{\mathbf{u}}$. Indeed, due to continuity of \mathcal{P} we obtain

$$\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \rightharpoonup \tilde{\mathbf{u}} \quad \text{in } L^2(0, T; L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

Thus with Proposition 3.7 and the compact embedding $L^{\frac{2\gamma}{\gamma+1}}(\mathbb{T}^N) \xhookrightarrow{c} W^{-1, 2}(\mathbb{T}^N)$

$$(3.21) \quad \mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \rightarrow \tilde{\mathbf{u}} \quad \text{in } L^2(0, T; W^{-1, 2}(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

Since

$$(3.22) \quad \operatorname{div} \tilde{\mathbf{u}}_\varepsilon \rightharpoonup 0 \quad \text{in } L^2(0, T; L^2(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

we have also that

$$(3.23) \quad \mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon \rightharpoonup \tilde{\mathbf{u}} \quad \text{in } L^2(0, T; W^{1, 2}(\mathbb{T}^N)) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

So combining (3.21) with (3.23) we conclude that

$$\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \cdot \mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon \rightharpoonup |\tilde{\mathbf{u}}|^2 \quad \text{in } L^1(Q) \quad \tilde{\mathbb{P}}\text{-a.s.}$$

Using Proposition 3.7 yields $\tilde{\mathbb{P}}$ -a.s.

$$\left| \int_Q (|\mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon|^2 - \mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \cdot \mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon) dx dt \right| \leq \|\tilde{\varrho}_\varepsilon - 1\|_{L^\infty(0,T;L^\gamma)} \|\tilde{\mathbf{u}}_\varepsilon\|_{L^2(0,T;L^s)}^2 \rightarrow 0,$$

where $s = \frac{2\gamma}{\gamma-1} < \frac{2N}{N-2}$. This implies $\|\mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon\|_2 \rightarrow \|\tilde{\mathbf{u}}\|_2$ and hence

$$\mathcal{P}_H \tilde{\mathbf{u}}_\varepsilon \rightarrow \tilde{\mathbf{u}} \quad \text{in } L^2(0,T;L^2(\mathbb{T}^N)).$$

Combining this with weak convergence in $L^2(0,T;W^{1,2}(\mathbb{T}^N))$ (recall Proposition 3.7) yields the claim. \square

In the following we aim to identify the limit in the gradient part of the convective term. To this end, we adopt the deterministic approach proposed by Lions and Masmoudi [18]. We introduce the dual space

$$W_{\text{div}}^{-l,2}(\mathbb{T}^N) \equiv \left[W_{\text{div}}^{l,2}(\mathbb{T}^N) \right]^*.$$

In particular, two elements of $W_{\text{div}}^{-l,2}(\mathbb{T}^N)$ are identical if their difference is a gradient.

Proposition 3.12. *For $l > \frac{N}{2}$ we have $\tilde{\mathbb{P}}$ -a.s.*

$$\text{div}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon \otimes \tilde{\mathbf{u}}_\varepsilon) \rightharpoonup \text{div}(\tilde{\mathbf{u}} \otimes \tilde{\mathbf{u}}) \quad \text{in } L^1(0,T;W_{\text{div}}^{-l,2}(\mathbb{T}^N)).$$

Proof. Following [18] we decompose

$$\begin{aligned} \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon &= \tilde{\mathbf{u}} + \mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) + \mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}), \\ \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon &= \tilde{\mathbf{u}} + \mathcal{P}_H(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) + \mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}). \end{aligned}$$

The claim follows once we can show that the following convergences hold true weakly in $L^1(0,T;W_{\text{div}}^{-l,2}(\mathbb{T}^N))$ $\tilde{\mathbb{P}}$ -a.s.:

$$(3.24) \quad \text{div} \left(\tilde{\mathbf{u}} \otimes \mathcal{P}_H(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

$$(3.25) \quad \text{div} \left(\tilde{\mathbf{u}} \otimes \mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

$$(3.26) \quad \text{div} \left(\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \tilde{\mathbf{u}} \right) \rightharpoonup 0,$$

$$(3.27) \quad \text{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \tilde{\mathbf{u}} \right) \rightharpoonup 0,$$

$$(3.28) \quad \text{div} \left(\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \mathcal{P}_H(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

$$(3.29) \quad \text{div} \left(\mathcal{P}_H(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

$$(3.30) \quad \text{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \mathcal{P}_H(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

$$(3.31) \quad \text{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \otimes \mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}) \right) \rightharpoonup 0,$$

The first four convergences follow from Proposition 3.7, (3.20) and the continuity of \mathcal{P}_H and \mathcal{Q} respectively. The convergences (3.28)-(3.30) are consequences of (3.17) and (3.18). In fact, the only critical part is (3.31). First, we need some improved space regularity. Similarly to [18], we use mollification by means of spatial convolution with a family of regularizing kernels with a parameter $\delta = \delta(\omega)$. As a matter of fact, thanks to the special geometry of the flat torus \mathbb{T}^N , the mollified functions can be taken as projections to a finite number of modes of the trigonometric basis $\{\exp(ikx)\}_{k \in \mathcal{Z}}$. In particular, the mollification commutes with all spatial derivatives as well as with the projections \mathcal{P}_h and \mathcal{Q} .

We take $\delta = \delta(\omega)$ so small that $\tilde{\mathbb{P}}$ -a.s.

$$(3.32) \quad \|(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta - \tilde{\varrho}_\varepsilon^\delta \tilde{\mathbf{u}}_\varepsilon^\delta\|_{L^2(L^{\frac{2\gamma}{\gamma+1}})} + \|(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta - \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon\|_{L^2(L^{\frac{2\gamma}{\gamma+1}})} \leq \delta,$$

$$(3.33) \quad \|(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta - \tilde{\mathbf{u}}_\varepsilon^\delta\|_{L^2(L^{\frac{2\gamma}{\gamma+1}})} + \|\tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}_\varepsilon\|_{L^2(L^{\frac{2N}{N-2}})} + \|\tilde{\mathbf{u}}^\delta - \tilde{\mathbf{u}}\|_{L^2(L^{\frac{2N}{N-2}})} \leq \delta,$$

uniformly in ε . We note that the norm $\|\tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}_\varepsilon\|_{L^2(L^{\frac{2N}{N-2}})}$ can be made uniformly small as a consequence of the gradient estimate (3.2).

As the mollification commutes with div and \mathcal{Q} , it is enough to show

$$(3.34) \quad \operatorname{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon^\delta \tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) \otimes \mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) \right) \rightharpoonup 0,$$

for fixed δ instead of (3.31). Second, we write

$$\mathcal{Q}(\tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) = \mathcal{Q}(\tilde{\varrho}_\varepsilon^\delta \tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) + \mathcal{Q}((1 - \tilde{\varrho}_\varepsilon^\delta) \tilde{\mathbf{u}}_\varepsilon^\delta).$$

By (3.17), the continuity of \mathcal{Q} and the boundedness of $\tilde{\mathbf{u}}_\varepsilon^\delta$ we know that

$$\mathcal{Q}((1 - \tilde{\varrho}_\varepsilon^\delta) \tilde{\mathbf{u}}_\varepsilon^\delta) \rightarrow 0 \quad \text{in } L^2(Q)$$

$\tilde{\mathbb{P}}$ -a.s. So (3.34) follows from

$$(3.35) \quad \operatorname{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon^\delta \tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) \otimes \mathcal{Q}(\tilde{\varrho}_\varepsilon^\delta \tilde{\mathbf{u}}_\varepsilon^\delta - \tilde{\mathbf{u}}^\delta) \right) \rightharpoonup 0,$$

in $L^1(0, T; W_{\operatorname{div}}^{-l, 2}(\mathbb{T}^N))$. As $\operatorname{div}(\mathcal{Q}\tilde{\mathbf{u}}^\delta \otimes \mathcal{Q}\tilde{\mathbf{u}}^\delta) = \frac{1}{2} \nabla |\mathcal{Q}\tilde{\mathbf{u}}^\delta|^2$, the convergence (3.35) is a consequence of

$$(3.36) \quad \operatorname{div} \left(\mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta \otimes \mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta \right) \rightharpoonup 0 \quad \text{in } L^1(0, T; W_{\operatorname{div}}^{-l, 2}(\mathbb{T}^N)),$$

thanks to (3.20) and (3.32). In order to show (3.36) (we need to introduce the function $\tilde{\Psi}_\varepsilon = \Delta^{-1} \operatorname{div}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)$ which satisfies $\nabla \tilde{\Psi}_\varepsilon = \mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)$). We have the system of equations

$$d(\varepsilon \tilde{\varphi}_\varepsilon) = -\nabla \tilde{\Psi}_\varepsilon \, dt, \quad d\nabla \tilde{\Psi}_\varepsilon = -\frac{\gamma}{\varepsilon} \nabla \tilde{\varphi}_\varepsilon \, dt + \tilde{\mathbf{F}}_\varepsilon \, dt + \mathcal{Q}\Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) d\tilde{W}_\varepsilon.$$

The right-hand-side only belongs to $W^{-l, 2}(\mathbb{T}^N)$. So we apply mollification and gain $\tilde{\Psi}_\varepsilon^\delta = \Delta^{-1} \operatorname{div}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta$ and $\nabla \tilde{\Psi}_\varepsilon^\delta = \mathcal{Q}(\tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta$. The system of equations for $\tilde{\varphi}_\varepsilon^\delta$ and $\tilde{\Psi}_\varepsilon^\delta$ reads as

$$(3.37) \quad d(\varepsilon \tilde{\varphi}_\varepsilon^\delta) = -\Delta \tilde{\Psi}_\varepsilon^\delta \, dt, \quad d\nabla \tilde{\Psi}_\varepsilon^\delta = -\frac{\gamma}{\varepsilon} \nabla \tilde{\varphi}_\varepsilon^\delta \, dt + \tilde{\mathbf{F}}_\varepsilon^\delta \, dt + \mathcal{Q}\Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta d\tilde{W}_\varepsilon.$$

We note that for the special choice, where the mollification is taken as the projection onto a finite number of Fourier modes, the system (3.37) reduces to a *finite number* of equations.

Now, we apply Itô's formula to the function

$$f(\varepsilon \tilde{\varphi}_\varepsilon^\delta, \nabla \tilde{\Psi}_\varepsilon^\delta) = \int_{\mathbb{T}^N} \varepsilon \tilde{\varphi}_\varepsilon^\delta \nabla \tilde{\Psi}_\varepsilon^\delta \cdot \varphi \, dx,$$

with $\varphi \in C_{\operatorname{div}}^\infty(\mathbb{T}^N)$ arbitrary and gain

$$\begin{aligned} & \int_{\mathbb{T}^N} \varepsilon \tilde{\varphi}_\varepsilon^\delta(t) \nabla \tilde{\Psi}_\varepsilon^\delta(t) \cdot \varphi \, dx \\ &= - \int_0^t \int_{\mathbb{T}^N} \Delta \tilde{\Psi}_\varepsilon^\delta \nabla \tilde{\Psi}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma - \gamma \int_0^t \int_{\mathbb{T}^N} \tilde{\varphi}_\varepsilon^\delta \nabla \tilde{\varphi}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma \\ &+ \varepsilon \int_0^t \int_{\mathbb{T}^N} \tilde{\varphi}_\varepsilon^\delta \tilde{\mathbf{F}}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma + \varepsilon \int_{\mathbb{T}^N} \int_0^t \tilde{\varphi}_\varepsilon^\delta \varphi \cdot \mathcal{Q}\Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta d\tilde{W}_\varepsilon \, dx. \end{aligned}$$

And we have

$$\int_0^t \int_{\mathbb{T}^N} \Delta \tilde{\Psi}_\varepsilon^\delta \nabla \tilde{\Psi}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma$$

$$\begin{aligned}
 &= \frac{1}{2} \int_0^t \int_{\mathbb{T}^N} \nabla |\nabla \tilde{\Psi}_\varepsilon^\delta|^2 \cdot \varphi \, dx \, d\sigma - \int_0^t \int_{\mathbb{T}^N} \nabla \tilde{\Psi}_\varepsilon^\delta \otimes \nabla \tilde{\Psi}_\varepsilon^\delta : \nabla \varphi \, dx \, d\sigma \\
 &= - \int_0^t \int_{\mathbb{T}^N} \nabla \tilde{\Psi}_\varepsilon^\delta \nabla \tilde{\Psi}_\varepsilon^\delta : \nabla \varphi \, dx \, d\sigma, \\
 &\int_0^t \int_{\mathbb{T}^N} \tilde{\varphi}_\varepsilon^\delta \nabla \tilde{\varphi}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma = \frac{1}{2} \int_0^t \int_{\mathbb{T}^N} \nabla |\tilde{\varphi}_\varepsilon^\delta|^2 \cdot \varphi \, dx \, d\sigma = 0,
 \end{aligned}$$

due to $\operatorname{div} \varphi = 0$. So we end up with

$$\begin{aligned}
 &\int_0^t \int_{\mathbb{T}^N} \nabla \tilde{\Psi}_\varepsilon^\delta \otimes \nabla \tilde{\Psi}_\varepsilon^\delta : \nabla \varphi \, dx \, d\sigma = -\varepsilon \int_{\mathbb{T}^N} \tilde{\varphi}_\varepsilon^\delta(t) \nabla \tilde{\Psi}_\varepsilon^\delta(t) \cdot \varphi \, dx \\
 &\quad + \varepsilon \int_0^t \int_{\mathbb{T}^N} \tilde{\varphi}_\varepsilon^\delta \tilde{\mathbf{F}}_\varepsilon^\delta \cdot \varphi \, dx \, d\sigma + \varepsilon \int_{\mathbb{T}^N} \int_0^t \tilde{\varphi}_\varepsilon^\delta \varphi \cdot \mathcal{Q}\Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)^\delta \, d\tilde{W}_\varepsilon \, dx.
 \end{aligned}$$

For fixed $\delta > 0$ the right-hand-side vanishes $\tilde{\mathbb{P}}$ -a.s. for $\varepsilon \rightarrow 0$ at least after taking a subsequence due to Corollary 3.10, Proposition 3.7 and the properties of the mollification. Finally we conclude with (3.36) which implies the last missing convergence (3.31) as explained above. \square

Now, we have all in hand to complete the proof of Theorem 3.8 which implies the proof of our main result, Theorem 2.8.

Proof of Theorem 3.8. The first part of the claim follows immediately from the fact that all \tilde{W}_ε are cylindrical Wiener processes due to Proposition 3.9. As a consequence, there exists a collection of mutually independent real-valued $(\tilde{\mathcal{F}}_t)$ -Wiener processes $(\tilde{\beta}_k)_{k \geq 1}$ such that $\tilde{W} = \sum_{k \geq 1} \tilde{\beta}_k e_k$.

In order to show that (1.2) is satisfied in the sense of Definition 2.4, let us take a divergence free test function $\varphi \in C_{\operatorname{div}}^\infty(\mathbb{T}^N)$ and consider the functionals M , N , N_k from Proposition 3.9. This way we only study the approximate equation (1.1b) projected by \mathcal{P}_H and the pressure term drops out. Having (3.14), (3.15) and (3.16) in hand, we intend to pass to the limit as $\varepsilon \rightarrow 0$ and to deduce

$$(3.38) \quad \tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\mathbf{u}}, \mathbf{r}_s \tilde{W}) [M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}})_{s,t}] = 0,$$

$$(3.39) \quad \tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\mathbf{u}}, \mathbf{r}_s \tilde{W}) \left[[M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}})^2]_{s,t} - N(1, \tilde{\mathbf{u}})_{s,t} \right] = 0,$$

$$(3.40) \quad \tilde{\mathbb{E}} h(\mathbf{r}_s \tilde{\mathbf{u}}, \mathbf{r}_s \tilde{W}) \left[[M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}}) \tilde{\beta}_k]_{s,t} - N_k(1, \tilde{\mathbf{u}})_{s,t} \right] = 0.$$

Note that the proof will then be complete. Indeed, (3.38), (3.39) and (3.40) imply that the process $M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}})$ is a $(\tilde{\mathcal{F}}_t)$ -martingale and its quadratic and cross variations satisfy, respectively,

$$\langle\langle M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}}) \rangle\rangle = N(1, \tilde{\mathbf{u}}), \quad \langle\langle M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}}), \tilde{\beta}_k \rangle\rangle = N_k(1, \tilde{\mathbf{u}}),$$

and consequently

$$\left\langle\left\langle M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}}) - \int_0^\cdot \langle \Phi(1, \tilde{\mathbf{u}}) \, d\tilde{W}, \varphi \rangle \right\rangle\right\rangle = 0$$

hence (1.2a) is satisfied in the sense required by Definition 2.4.

Let us now verify (3.38), (3.39) and (3.40). First of all we observe that

$$M(\tilde{\varrho}_\varepsilon, \tilde{\mathbf{u}}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon)_t \rightarrow M(1, \tilde{\mathbf{u}}, \tilde{\mathbf{u}})_t \quad \text{a.s.}$$

due to Proposition 3.7, Proposition 3.12 and (3.21). Application of the Vitali convergence theorem together with the uniform estimates (3.4), (3.8) and (3.9) justifies the passage

to the limit in (3.14) and (3.38) follows. The same argument implies the passage to the limit in the part of (3.15) and (3.16) involving M .

Finally, we comment on the passage to the limit in the terms coming from the stochastic integral, i.e. N and N_k . The convergence in (3.16) being easier, let us only focus on (3.15) in detail. As the first step we note that the convergence

$$\sum_{k \geq 1} \langle \mathbf{g}_k(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon), \boldsymbol{\varphi} \rangle^2 \rightarrow \sum_{k \geq 1} \langle \mathbf{g}_k(1, \tilde{\mathbf{u}}), \boldsymbol{\varphi} \rangle^2 \quad \tilde{\mathbb{P}} \otimes \mathcal{L}\text{-a.e.}$$

follows once we show that

$$(3.41) \quad \langle \Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \cdot, \boldsymbol{\varphi} \rangle \rightarrow \langle \Phi(1, \tilde{\mathbf{u}}) \cdot, \boldsymbol{\varphi} \rangle \quad \text{in } L_2(\mathcal{M}; \mathbb{R}) \quad \tilde{\mathbb{P}} \otimes \mathcal{L}\text{-a.e.}$$

To this end, we write

$$\begin{aligned} & \left\| \langle \Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \cdot, \boldsymbol{\varphi} \rangle - \langle \Phi(1, \tilde{\mathbf{u}}) \cdot, \boldsymbol{\varphi} \rangle \right\|_{L_2(\mathcal{M}; \mathbb{R})} \\ & \leq \left(\sum_{k \geq 1} |\langle \mathbf{h}_k(\tilde{\varrho}_\varepsilon) - \mathbf{h}_k(1), \boldsymbol{\varphi} \rangle|^2 \right)^{\frac{1}{2}} + \left(\sum_{k \geq 1} |\alpha_k|^2 |\langle \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon - \tilde{\mathbf{u}}, \boldsymbol{\varphi} \rangle|^2 \right)^{\frac{1}{2}} \\ & = I_1 + I_2. \end{aligned}$$

For I_2 we use (2.2) together with (3.21) to obtain $I_2 \rightarrow 0$ for a.e. (ω, t) . For I_1 we apply the Minkowski integral inequality, the mean value theorem, (2.3) and (2.4) to obtain

$$\begin{aligned} I_1 & \leq C \left(\sum_{k \geq 1} \|\mathbf{h}_k(\tilde{\varrho}_\varepsilon) - \mathbf{h}_k(1)\|_{L_x^1}^2 \right)^{\frac{1}{2}} \leq C \int_{\mathbb{T}^N} \left(\sum_{k \geq 1} |\mathbf{h}_k(\tilde{\varrho}_\varepsilon) - \mathbf{h}_k(1)|^2 \right)^{\frac{1}{2}} dx \\ & \leq C \int_{\mathbb{T}^N} \left(1 + \tilde{\varrho}_\varepsilon^{\frac{\gamma-1}{2}} \right) |\tilde{\varrho}_\varepsilon - 1| dx \leq C \left[\int_{\mathbb{T}^N} \left(1 + \tilde{\varrho}_\varepsilon^{\frac{\gamma-1}{2}} \right)^p dx \right]^{\frac{1}{p}} \left[\int_{\mathbb{T}^N} |\tilde{\varrho}_\varepsilon - 1|^q dx \right]^{\frac{1}{q}} \end{aligned}$$

where the conjugate exponents $p, q \in (1, \infty)$ are chosen in such a way that

$$p \frac{\gamma-1}{2} < \gamma + 1 \quad \text{and} \quad q < \gamma.$$

Therefore, using (3.7), (3.17) we deduce

$$\tilde{\mathbb{E}} \int_0^T I_1 dt \rightarrow 0.$$

and so for a subsequence $I \rightarrow 0$ for a.e. (ω, t) and (3.41) follows. Besides, since, for all $p \geq 2$,

$$\begin{aligned} & \tilde{\mathbb{E}} \int_s^t \left\| \langle \Phi(\tilde{\varrho}_\varepsilon, \tilde{\varrho}_\varepsilon \tilde{\mathbf{u}}_\varepsilon) \cdot, \boldsymbol{\varphi} \rangle \right\|_{L_2(\mathcal{M}; \mathbb{R})}^p dr \\ & \leq C \tilde{\mathbb{E}} \int_s^t \|\tilde{\varrho}_\varepsilon\|_{L^2}^{\frac{p}{2}} \left(1 + \|\tilde{\varrho}_\varepsilon\|_{L^\gamma}^\gamma + \|\sqrt{\tilde{\varrho}_\varepsilon} \tilde{\mathbf{u}}_\varepsilon\|_{L^2}^2 \right)^{\frac{p}{2}} dr \\ & \leq C \left(1 + \tilde{\mathbb{E}} \sup_{0 \leq t \leq T} \|\tilde{\varrho}_\varepsilon\|_{L^\gamma}^p + \tilde{\mathbb{E}} \sup_{0 \leq t \leq T} \|\sqrt{\tilde{\varrho}_\varepsilon} \tilde{\mathbf{u}}_\varepsilon\|_{L^2}^{2p} \right) \leq C \end{aligned}$$

due to (3.3), (3.7), we obtain the convergence in (3.15) and therefore $\tilde{\mathbf{u}}$ solves (1.2). It follows immediately from our construction that for all $p \in [1, \infty)$

$$\tilde{\mathbf{u}} \in L^p(\tilde{\Omega}; L^2(0, T; W_{\text{div}}^{1,2}(\mathbb{T}^N))).$$

Besides, since we have (due Proposition 3.7 and (3.17))

$$\sqrt{\tilde{\varrho}_\varepsilon} \tilde{\mathbf{u}}_\varepsilon \rightharpoonup \tilde{\mathbf{u}} \quad \text{in } L^1(\Omega; L^1(Q))$$

lower semi-continuity of the functional

$$\tilde{\mathbf{w}} \mapsto \tilde{\mathbb{E}} \left[\sup_{t \in (0, T)} \int_{\mathbb{T}^N} |\tilde{\mathbf{w}}|^2 dx \right]^{\frac{p}{2}}$$

yields $\tilde{\mathbf{u}} \in L^p(\tilde{\Omega}; L^\infty(0, T; L^2(\mathbb{T}^N)))$ on account of Corollary 3.10. The usual argument about the fractional time derivative (in the distributional sense) implies

$$\tilde{\mathbf{u}} \in L^p(\tilde{\Omega}; C_w([0, T]; L^2_{\text{div}}(\mathbb{T}^N)))$$

and the proof is complete. \square

4. PROOF OF THEOREM 2.9

In order to complete the proof of Theorem 2.9, we make use of Proposition A.4 which is a generalization of the Gyöngy-Krylov characterization of convergence in probability introduced in [13] adapted to the case of quasi-Polish spaces. It applies to situations when pathwise uniqueness and existence of a martingale solution are valid and allows to establish existence of a pathwise solution. We recall that in the case of $N = 2$ pathwise uniqueness for (1.2) is known (cf. Theorem 2.5).

We consider the collection of joint laws of

$$(\varrho_n, \mathbf{u}_n, \mathcal{P}(\varrho_n \mathbf{u}_n), \varrho_m, \mathbf{u}_m, \mathcal{P}(\varrho_m \mathbf{u}_m)) \quad \text{on} \quad \mathcal{X}_\varrho \times \mathcal{X}_{\mathbf{u}} \times \mathcal{X}_{\varrho \mathbf{u}} \times \mathcal{X}_\varrho \times \mathcal{X}_{\mathbf{u}} \times \mathcal{X}_{\varrho \mathbf{u}},$$

denoted by $\mu^{n, m}$. For this purpose we define the extended path space

$$\mathcal{X}^J = \mathcal{X}_\varrho \times \mathcal{X}_{\mathbf{u}} \times \mathcal{X}_{\varrho \mathbf{u}} \times \mathcal{X}_\varrho \times \mathcal{X}_{\mathbf{u}} \times \mathcal{X}_{\varrho \mathbf{u}} \times \mathcal{X}_W$$

As above, denote by μ_W the law of W and set $\nu^{n, m}$ to be the joint law of

$$(\varrho_n, \mathbf{u}_n, \mathcal{P}(\varrho_n \mathbf{u}_n), \varrho_m, \mathbf{u}_m, \mathcal{P}(\varrho_m \mathbf{u}_m), W) \quad \text{on} \quad \mathcal{X}^J.$$

Similarly to Corollary 3.6 the following fact holds true. The proof is nearly identical and so will be left to the reader.

Proposition 4.1. *The collection $\{\nu^{n, m}; n, m \in \mathbb{N}\}$ is tight on \mathcal{X}^J .*

Let us take any subsequence $\{\nu^{n_k, m_k}; k \in \mathbb{N}\}$. By the Jakubowski-Skorokhod theorem, Theorem A.2, we infer (for a further subsequence but without loss of generality we keep the same notation) the existence a probability space $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$ with a sequence of random variables

$$(\hat{\varrho}_{n_k}, \hat{\mathbf{u}}_{n_k}, \hat{\mathbf{q}}_{n_k}, \check{\varrho}_{m_k}, \check{\mathbf{u}}_{m_k}, \check{\mathbf{q}}_{m_k}, \bar{W}_k), \quad k \in \mathbb{N},$$

converging almost surely in \mathcal{X}^J to a random variable

$$(\hat{\varrho}, \hat{\mathbf{u}}, \hat{\mathbf{q}}, \check{\varrho}, \check{\mathbf{u}}, \check{\mathbf{q}}, \bar{W})$$

and

$$\tilde{\mathbb{P}}((\hat{\varrho}_{n_k}, \hat{\mathbf{u}}_{n_k}, \hat{\mathbf{q}}_{n_k}, \check{\varrho}_{m_k}, \check{\mathbf{u}}_{m_k}, \check{\mathbf{q}}_{m_k}, \bar{W}_k) \in \cdot) = \nu^{n_k, m_k}(\cdot).$$

Observe that in particular, μ^{n_k, m_k} converges weakly to a measure μ defined by

$$\mu(\cdot) = \tilde{\mathbb{P}}((\hat{\varrho}, \hat{\mathbf{u}}, \hat{\mathbf{q}}, \check{\varrho}, \check{\mathbf{u}}, \check{\mathbf{q}}) \in \cdot).$$

As the next step, we should recall the technique established in Subsection 3.4. Analogously, it can be applied to both

$$(\hat{\varrho}_{n_k}, \hat{\mathbf{u}}_{n_k}, \hat{\mathbf{q}}_{n_k}, \bar{W}_k), \quad (\hat{\varrho}, \hat{\mathbf{u}}, \hat{\mathbf{q}}, \bar{W})$$

and

$$(\check{\varrho}_{m_k}, \check{\mathbf{u}}_{m_k}, \check{\mathbf{q}}_{m_k}, \bar{W}_k), \quad (\check{\varrho}, \check{\mathbf{u}}, \check{\mathbf{q}}, \bar{W})$$

in order to show that $(\hat{\mathbf{u}}, \bar{W})$ and $(\check{\mathbf{u}}, \bar{W})$ are weak martingale solutions to (1.2) defined on the same stochastic basis $(\bar{\Omega}, \bar{\mathcal{F}}, (\bar{\mathcal{F}}_t), \bar{\mathbb{P}})$, where $(\bar{\mathcal{F}}_t)$ is the $\bar{\mathbb{P}}$ -augmented canonical filtration of $(\hat{\mathbf{u}}, \check{\mathbf{u}}, \bar{W})$. Besides, we obtain that

$$\hat{\varrho} = \check{\varrho}, \quad \hat{\mathbf{q}} = \check{\mathbf{q}}, \quad \check{\mathbf{q}} = \hat{\mathbf{q}} \quad \bar{\mathbb{P}}\text{-a.s.}$$

In order to verify the condition (A.1) from Proposition A.4 we employ the pathwise uniqueness result for (1.2) in two dimensions, cf. Theorem 2.5. Indeed, it follows from our assumptions on the approximate initial laws Λ_ε that $\hat{\mathbf{u}}(0) = \check{\mathbf{u}}(0) = \mathbf{u}_0$ $\bar{\mathbb{P}}$ -a.s., therefore according to Theorem 2.7 the solutions $\hat{\mathbf{u}}$ and $\check{\mathbf{u}}$ coincide $\bar{\mathbb{P}}$ -a.s. and

$$\begin{aligned} & \mu\left((\varrho_1, \mathbf{u}_1, \mathbf{q}_1, \varrho_2, \mathbf{u}_2, \mathbf{q}_2); (\varrho_1, \mathbf{u}_1, \mathbf{q}_1) = (\varrho_2, \mathbf{u}_2, \mathbf{q}_2)\right) \\ &= \bar{\mathbb{P}}\left((\hat{\varrho}, \hat{\mathbf{u}}, \hat{\mathbf{q}}) = (\check{\varrho}, \check{\mathbf{u}}, \check{\mathbf{q}})\right) = \bar{\mathbb{P}}(\hat{\mathbf{u}} = \check{\mathbf{u}}) = 1. \end{aligned}$$

Now, we have all in hand to apply Proposition A.4. It implies that the original sequence $(\varrho_\varepsilon, \mathbf{u}_\varepsilon, \mathcal{P}(\varrho_\varepsilon \mathbf{u}_\varepsilon))$ defined on the initial probability space $(\Omega, \mathcal{F}, \mathbb{P})$ converges in probability in the topology of $\mathcal{X}_\varrho \times \mathcal{X}_\mathbf{u} \times \mathcal{X}_{\varrho\mathbf{u}}$ to a random variable $(\varrho, \mathbf{u}, \mathbf{q})$. Without loss of generality, we assume that the convergence is almost sure and again by the method from Subsection 3.4 we finally deduce that \mathbf{u} is a pathwise weak solution to (1.2). Actually, identification of the limit is more straightforward here since in this case all the work is done for the initial setting and only one fixed driving Wiener process W is considered. The proof of Theorem 2.9 is complete.

APPENDIX A. QUASI-POLISH SPACES

The so-called quasi-Polish spaces are topological spaces that are not necessarily metrizable but nevertheless they enjoy several important properties of Polish spaces. Let us recall their definition introduced in [14].

Definition A.1. *Let (X, τ) be a topological space such that there exists a countable family*

$$\{f_n : X \rightarrow [-1, 1]; n \in \mathbb{N}\}$$

of continuous functions that separate points of X .

Among the properties of quasi-Polish spaces used in the main body of this paper belongs the following Jakubowski-Skorokhod representation theorem, see [14, Theorem 2].

Theorem A.2. *Let (X, τ) be a quasi-Polish space and let \mathcal{S} be the σ -field generated by $\{f_n; n \in \mathbb{N}\}$. If $\{\mu_n; n \in \mathbb{N}\}$ is a tight sequence of probability measures on (X, \mathcal{S}) , then there exists a subsequence (n_k) , a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with X -valued Borel measurable random variables $\{\xi_k; k \in \mathbb{N}\}$ and ξ such that μ_{n_k} is the law of ξ_k and ξ_k converges to ξ in X a.s. Moreover, the law of ξ is a Radon measure.*

Next, we need to adapt the Gyöngy-Krylov characterization of convergence in probability introduced in [13] to the setting of quasi-Polish spaces. Recall that the original argument for the case of Polish spaces follows from the following simple observation made in [13, Lemma 1.1].

Lemma A.3. *Let X be a Polish space equipped with the Borel σ -algebra. A sequence of X -valued random variables $\{Y_n; n \in \mathbb{N}\}$ converges in probability if and only if for every subsequence of joint laws, $\{\mu_{n_k, m_k}; k \in \mathbb{N}\}$, there exists a further subsequence which converges weakly to a probability measure μ such that*

$$\mu((x, y) \in X \times X; x = y) = 1.$$

In view of our application in Subsection 4, we are interested in the sufficiency of the above condition.

Proposition A.4. *Let (X, τ) be a quasi-Polish space. Let $\{Y_n; n \in \mathbb{N}\}$ be a sequence of X -valued random variables. Assume that for every subsequence of their joint laws $\{\mu_{n_k, m_k}; k \in \mathbb{N}\}$ there exists a further subsequence which converges weakly to a probability measure μ such that*

$$(A.1) \quad \mu((x, y) \in X \times X; x = y) = 1.$$

Then there exists a subsequence $\{Y_{n_l}; l \in \mathbb{N}\}$ which converges a.s.

Proof. Let \tilde{f} be the one-to-one and continuous mapping defined by

$$\begin{aligned} \tilde{f} : X &\rightarrow [-1, 1]^{\mathbb{N}} \\ x &\mapsto \{f_n(x); n \in \mathbb{N}\}, \end{aligned}$$

where f_n were given by Definition A.1. Since due to assumption

$$(Y_{n_k}, Y_{m_k}) \xrightarrow{d} (Y, Y) \quad \text{in } X \times X$$

for every $(n_k), (m_k)$ and some Y with the law μ , we deduce from the continuous mapping theorem that

$$(\tilde{f}(Y_{n_k}), \tilde{f}(Y_{m_k})) \xrightarrow{d} (\tilde{f}(Y), \tilde{f}(Y)) \quad \text{in } [-1, 1]^{\mathbb{N}} \times [-1, 1]^{\mathbb{N}}$$

for every $(n_k), (m_k)$. Since $[-1, 1]^{\mathbb{N}} \times [-1, 1]^{\mathbb{N}}$ is a Polish space, Lemma A.3 applies to the sequence $\{\tilde{f}(Y_n); n \in \mathbb{N}\}$ and the convergence in probability follows. Consequently, there exists a subsequence $\{\tilde{f}(Y_{n_l}); n \in \mathbb{N}\}$ which converges a.s. and it only remains to prove that Y_{n_l} converges to Y a.s. To this end, we proceed by contradiction: Assume that Y_{n_l} does not converge a.s. to Y . Then there exists a set of positive probability $\Omega^* \subset \Omega$ such that for all $\omega \in \Omega^*$ there exists a neighborhood $\mathcal{N}(\omega)$ of $Y(\omega)$ and for every $l_0 \in \mathbb{N}$ there exists $l \geq l_0$ such that $Y_{n_l}(\omega) \notin \mathcal{N}(\omega)$. However, as the sequence $\{f_n; n \in \mathbb{N}\}$ separates points of X , there exists $n \in \mathbb{N}$ such that $f_n(Y_{n_l}(\omega)) \neq f_n(Y(\omega))$ and as a consequence there exists a neighborhood $\mathcal{V}(\omega)$ of $\tilde{f}(Y(\omega))$ such that $\tilde{f}(Y_{n_l}(\omega)) \notin \mathcal{V}(\omega)$. This contradicts the a.s. convergence of $\tilde{f}(Y_{n_l})$ and completes the proof. \square

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Conflict of Interest

The authors declare that they have no conflict of interest.

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