

Rigorous results in space-periodic two-dimensional turbulence

Sergei Kuksin^a and Armen Shirikyan^{b,c}

^a CNRS, Institut de Mathématiques de Jussieu–Paris Rive Gauche, UMR 7586,
Université Paris Diderot, Sorbonne Paris Cité, F-75013, Paris, France
E-mail: Sergei.Kuksin@imj-prg.fr

^b Département de Mathématiques, Université de Cergy–Pontoise, CNRS UMR8088
2 avenue Adolphe Chauvin, 95302 Cergy–Pontoise Cedex, France
E-mail: Armen.Shirikyan@u-cergy.fr

^c Centre de Recherches Mathématiques, CNRS UMI3457, Université de Montréal
Montréal, QC, H3C 3J7, Canada

Abstract

We survey the recent advance in the rigorous qualitative theory of the 2d stochastic Navier–Stokes system that are relevant to the description of turbulence in two-dimensional fluids. After discussing briefly the initial-boundary value problem and the associated Markov process, we formulate results on the existence, uniqueness and mixing of a stationary measure. We next turn to various consequences of these properties: strong law of large numbers, central limit theorem, and random attractors related to a unique stationary measure. We also discuss the Donsker–Varadhan and Freidlin–Wentzell type large deviations, as well as the inviscid limit and asymptotic results in 3d thin domains. We conclude with some open problems.

AMS subject classifications: 35Q30, 35R60, 37A25, 37L55, 60F05, 60F10, 60H15, 76D05,

Keywords: 2d Navier–Stokes system, stationary measure, mixing, strong law of large numbers, central limit theorem, random attractors, large deviations, inviscid limit

Contents

0	Introduction	2
1	Equations and random forces	3
1.1	Cauchy problem	3
1.2	Markov process and a priori estimates	6

2	Mixing	8
2.1	Existence of a stationary measure	8
2.2	Uniqueness and exponential stability	9
3	Consequences of mixing	12
3.1	Ergodic theorems	12
3.2	Random attractors	13
3.3	Dependence on parameters and stability	15
4	Large deviations	15
4.1	Donsker–Varadhan type large deviations	16
4.2	Vanishing noise limit	18
5	Inviscid limit	19
5.1	Properties of μ^ν , independent from ν	19
5.2	Inviscid limit	20
6	3d Navier–Stokes system in thin domains	22
7	Open problems	24
	Bibliography	26

0 Introduction

The main subject of this article is the two-dimensional Navier–Stokes system in \mathbb{R}^2 subject to periodic boundary conditions, perturbed by a random force. We thus consider the equations

$$\partial_t u + \langle u, \nabla \rangle u - \nu \Delta u + \nabla p = f(t, x), \quad \operatorname{div} u = 0, \quad (0.1)$$

where $u = (u_1, u_2)$ and p are the velocity field and pressure of a fluid, $\nu > 0$ is the kinematic viscosity, f is an external (random) force, and $\langle u, \nabla \rangle = u_1 \partial_1 + u_2 \partial_2$. All the functions are assumed to be 2π -periodic in the spacial variables x_1 and x_2 . Equations (0.1) are supplemented with the initial condition

$$u(0, x) = u_0(x), \quad (0.2)$$

where u_0 is a given divergence-free vector field that is 2π -periodic and locally square-integrable. Our aim is to review rigorous results on the qualitative behaviour of solutions for (0.1) as $t \rightarrow \infty$ and/or $\nu \rightarrow 0$. These two limits are important for the mathematical description of the space-periodic 2d turbulence (for a physical treatment of this topic see e.g. [BV12]). In our review we avoid detailed discussion, related to the relevance of the 2d Navier–Stokes system (0.1) for physics, referring the reader to [Bat82, Fri95, Gal02, BV12]. But we mention that a number of equations, closely related to the 2d Navier–Stokes system, are

used in meteorology and oceanography, so the methods, developed for Navier–Stokes equations, can be applied in the mathematical theory of climate and the statistical description of the ocean. A result, presented in Section 6, gives a very basic explanation for the relevance of the 2d models for statistical description of the 3d phenomena.

Most of the results, discussed in our work, were obtained in this century. Complete proofs and discussion of many of them can be found in [KS12]. New material, not treated in that book, includes the theory of large deviations and recent progress concerning the mixing in Eq. (0.1).

Acknowledgement. This research was supported by the *Agence Nationale de la Recherche* through the grants ANR-10-BLAN 0102 and ANR-17-CE40-0006-02. The research of AS was carried out within the MME-DII Center of Excellence (ANR-11-LABX-0023-01) and supported by the *Initiative d’excellence Paris-Seine* and by the CNRS PICS *Fluctuation theorems in stochastic systems*.

1 Equations and random forces

1.1 Cauchy problem

As was mentioned in the Introduction, we consider Eqs. (0.1) with periodic boundary conditions. Before describing the class of random forces f we deal with, let us recall a general result on the existence, uniqueness, and regularity of solutions to the deterministic Cauchy problem (0.1), (0.2). We begin with the definition of a solution on an arbitrary time interval $J_T = [0, T]$. The definitions of all the functional spaces used here and henceforth can be found in the list of frequently used notations at the end of the paper.

Definition 1.1. Let $f(t, x)$ be the time derivative of a piecewise continuous function¹ $g : J_T \rightarrow L^2(\mathbb{T}^2, \mathbb{R}^2)$, vanishing at $t = 0$: $f = \partial_t g(t, x)$. A function $u(t, x)$ defined on $J_T \times \mathbb{T}^2$ is called a *weak solution* for (0.1) if it belongs to the space

$$\mathcal{X}_T = C(J_T, L^2_\sigma) \cap L^2(J_T, H^1(\mathbb{T}^2, \mathbb{R}^2) \cap L^2_\sigma),$$

and satisfies the relation

$$(u(t), \varphi) + \int_0^t (\langle u, \nabla \rangle u - \nu \Delta u, \varphi) ds = (u(0), \varphi) + (g(t), \varphi), \quad t \in J_T, \quad (1.1)$$

where φ is an arbitrary divergence-free smooth vector field on \mathbb{T}^2 , and the term with the Laplacian under the integral is understood in the weak sense: $(\Delta u, \varphi) = -(\nabla u, \nabla \varphi)$. In what follows, when discussing Eqs. (0.1), we often say *solution* rather than *weak solution*.

¹Here and in similar situations below, it means that g has at most countably many points of discontinuity, where it has left and right limits. Traditionally we normalise such functions to be right-continuous everywhere. In particular, $g(0)$ is well defined.

Note that the Navier–Stokes system contains three unknown functions, u_1 , u_2 , and p , whereas the definition of a solution specifies only the velocity field u . This is due to the fact that, once $u(t, x)$ satisfying (1.1) is constructed, the Leray decomposition (see Theorem 1.5 in [Tem79, Chapter 1]) may be used to find a unique (up to an additive function of t) distribution $p(t, x)$ such that the first equation in (0.1) holds in a weak sense. A proof of the following theorem is essentially contained in [Tem79, Chapter 3].

Theorem 1.2. *Let $T > 0$ and let f be the sum of a square-integrable function $h(x)$ and the time derivative of a piecewise continuous function with range in $H^1(\mathbb{T}^2, \mathbb{R}^2)$. Then, for any $u_0 \in L^2_\sigma$, Eq. (0.1) has a unique weak solution $u \in \mathcal{X}_T$ satisfying the initial condition (0.2).*

If the function $g(t, x)$ is such that its mean value (in x) vanishes identically in t , then the mean value of the solution $u(t, x)$ is time-independent. Below we always assume that the mean values of the forces we apply to the Navier–Stokes system and of its solutions which we consider, both vanish identically in time.

Theorem 1.2 allows one to construct a unique solution of Eq. (0.1) for the three important classes of random forces f , specified below. Namely, let H be the space of divergence free square-integrable vector fields on \mathbb{T}^2 with zero mean value. We shall deal with random forces of the form

$$f(t, x) = h(x) + \eta(t, x), \quad (1.2)$$

where $h \in H$ is a deterministic function, and η is one of the following three random processes:

Spatially regular white noise. Let \mathbb{Z}_*^2 be the set of non-zero integer vectors $j = (j_1, j_2)$ and let $\{e_j, j \in \mathbb{Z}_*^2\}$ be a trigonometric basis in H defined by

$$e_j(x) = \frac{j^\perp}{\sqrt{2}\pi|j|} \begin{cases} \cos\langle j, x \rangle & \text{if } j_1 > 0 \text{ and if } j_1 = 0, j_2 > 0, \\ \sin\langle j, x \rangle & \text{if } j_1 < 0 \text{ and if } j_1 = 0, j_2 < 0, \end{cases} \quad (1.3)$$

where $j^\perp = (-j_2, j_1)$. Let us fix numbers $\{b_j, j \in \mathbb{Z}_*^2\}$ such that

$$\mathfrak{B}_1 < \infty,$$

where for $k \geq 0$ we denote

$$\mathfrak{B}_k = \sum_{j \in \mathbb{Z}_*^2} b_j^2 |j|^{2k} \leq \infty, \quad (1.4)$$

and define

$$\eta(t, x) = \frac{\partial}{\partial t} \zeta(t, x), \quad \zeta(t, x) = \sum_{j \in \mathbb{Z}_*^2} b_j \beta_j(t) e_j(x), \quad (1.5)$$

where $\{\beta_j, j \in \mathbb{Z}_*^2\}$ is a family of independent standard Brownian motions. Then $\eta(t, x) = \sum b_j \eta_j(t) e_j(x)$, where $\{\eta_j = \dot{\beta}_j, j \in \mathbb{Z}_*^2\}$, are standard independent

white noises. It follows from the Doob–Kolmogorov inequality (see Theorem 3.8 in [KS91, Chapter 1]) that, with probability 1, the series in (1.5) converges in H^1 uniformly in $t \in [0, T]$ for any $T < \infty$, so that $\eta(t)$ is the time derivative of a continuous vector function with range in the space $V = H \cap H^1(\mathbb{T}^2, \mathbb{R}^2)$.

Random kicks. Let $\{\eta^k\}$ be a sequence of i.i.d. random variables (the kicks) with range in V . Define

$$\eta(t, x) = \frac{\partial}{\partial t} \zeta(t, x), \quad \zeta(t, x) = \sum_{k=1}^{\infty} \eta^k(x) \theta(t - k), \quad (1.6)$$

where $\theta(t) = 0$ for $t < 0$ and $\theta(t) = 1$ for $t \geq 0$. Since η has jumps only at positive integers, the trajectories of η are the time derivative of piecewise continuous functions.

Piecewise independent process. Let $\{\eta^k\}$ be a sequence of i.i.d. random variables in $L^2(J_1, H)$. We define a random process of the form

$$\eta(t, x) = \sum_{k=1}^{\infty} I_{[k-1, k)}(t) \eta^k(t - k + 1, x), \quad (1.7)$$

where $I_{[k-1, k)}$ is the indicator function of the interval.

Example 1.3. Let us take H -valued processes $\{\eta^k(t, \cdot), 0 \leq t \leq 1\}$ of the form

$$\eta^k(t, x) = \sum_{j \in \mathbb{Z}_*^2} b_j \eta_j^k(t) e_j(x), \quad \mathfrak{B}_1 < \infty,$$

where $\{\eta_j^k : k \geq 1, j \in \mathbb{Z}_*^2\}$ are real-valued independent random processes, distributed as a fixed process $\tilde{\eta} : [0, 1] \rightarrow \mathbb{R}$. Taking for $\tilde{\eta}$ a random series with respect to the Haar basis of the space $L_2(0, 1)$ (see Section 21 in [Lam96]), we arrive at a random process η as in (1.7) that has the form

$$\eta(t, x) = \sum_{j \in \mathbb{Z}_*^2} b_j \eta_j(t) e_j(x). \quad (1.8)$$

Here $\{\eta_j\}$ are independent random processes, distributed as the process

$$\eta_0(t) = c \sum_{l=1}^{\infty} \xi_l I_{[l-1, l)}(t) + \sum_{N=0}^{\infty} \sum_{l=0}^{\infty} c_n \xi_l^n H_l^n(t), \quad (1.9)$$

where c and c_n are real constants, ξ_l and ξ_l^n are i.i.d. real-valued random variables with a law λ , and H_l^n are the Haar functions:

$$H_l^n(t) = \begin{cases} 0 & \text{if } t < l2^{-n} \text{ or } t \geq (l+1)2^{-n}, \\ 1 & \text{if } l2^{-n} \leq t < (l+1/2)2^{-n}, \\ -1 & \text{if } (l+1/2)2^{-n} \leq t < (l+1)2^{-n}. \end{cases}$$

Thus, $\eta_0(t)$ is a random wavelet series, and $\eta(t, x)$ is an H -valued process whose expansion in the trigonometric basis $\{e_j\}$ has independent random wavelet coefficients. Notice that a.s. the set of discontinuities of η is the family of dyadic numbers and, hence, is dense on the positive half-line. Besides, all trajectories of η are continuous at non-dyadic points and right-continuous everywhere.

If $c = 1$, $c_n = 2^{n/2}$, and λ is the centred normal law with unit variance, then $\eta_0(t)$ is white noise (see [Lam96]). If $c_n \ll 2^{n/2}$, then $\eta_0(t)$ is a *red noise*. In particular, if $|c_n| \leq Cn^{-q}$ for some $q > 1$, and λ has a bounded support, then the red noise $\eta_0(t)$ is bounded uniformly in t and ω .

For reasons of space, we shall usually state the results for the case of spatially regular white noise. However, suitable reformulations of most of them remain valid in the two other cases.

1.2 Markov process and a priori estimates

The family of solutions, corresponding to the three classes of random forces considered in the foregoing subsection, give rise to Markov processes in the space H . This allows us to apply to the study of the Navier–Stokes system (0.1) with such random forces well-developed probabilistic methods. We start with Eqs. (0.1), (1.5) (recalling that always $h \in H$ and $\mathfrak{B}_1 < \infty$), and for a random initial data $u_0 = u_0^\omega(x)$, independent from the force f , denote by $u(t; u_0)$ a solution of (0.1), (0.2), (1.5). For a non-random $u_0 = v \in H$ denote

$$P_t(v, \cdot) = \mathcal{D}(u(t; v)).$$

This is a measure in H , depending on (t, v) in a measurable way, and satisfying the Kolmogorov–Chapman relation. So P_t is the transition function of a Markov process in H . The latter is the Markov process generated by solutions of Eqs. (0.1), (1.5). It defines the *Markov semigroups* on functions and on measures by the relations

$$\mathfrak{P}_t : C_b(H) \rightarrow C_b(H), \quad (\mathfrak{P}_t g)(u) = \int_H P_t(u, dz) g(z), \quad (1.10)$$

$$\mathfrak{P}_t^* : \mathcal{P}(H) \rightarrow \mathcal{P}(H), \quad (\mathfrak{P}_t^* \mu)(\Gamma) = \int_H P_t(z, \Gamma) \mu(dz). \quad (1.11)$$

The two semigroups are instrumental to study the equation since the former defines the evolution of mean values of observables: for any $g \in C_b(H)$ and any $v \in H$, we have

$$\mathbb{E} g(u(t; v)) = (\mathfrak{P}_t g)(v), \quad t \geq 0. \quad (1.12)$$

On the other hand, the latter defines the evolution of the laws since

$$\mathcal{D}(u(t; u_0)) = \mathfrak{P}_t^* \mu, \quad t \geq 0, \quad (1.13)$$

if u_0 is a random variable independent from f , and its law equals μ .

The white in time structure of the noise allows not only to prove the Markovian character of evolution, but also to derive a priori estimates by an application of Ito's formula. The theorem below summarises some of them. For any integer $k \geq 0$, we set

$$\mathcal{E}_u(k, t) = t^k \|u(t)\|_k^2 + \int_0^t s^k \|u(s)\|_{k+1}^2 ds, \quad t \geq 0.$$

In the case $k = 0$, we write $\mathcal{E}_u(t)$. For $k \in \mathbb{N}$ we denote $H^k = H \cap H^k(\mathbb{T}^2, \mathbb{R}^2)$ (so $H^1 = V$).

Theorem 1.4. *Consider Eqs. (0.1), (1.5). The following properties hold for any $\nu > 0$ and any H -valued random variable u_0 , independent from ζ .*

Energy and enstrophy balances. *If $\mathbb{E}|u_0|_2^2 < \infty$, then*

$$\mathbb{E}|u(t)|_2^2 + 2\nu \mathbb{E} \int_0^t |\nabla u(s)|_2^2 ds = \mathbb{E}|u_0|_2^2 + \mathfrak{B}_0 t + 2 \mathbb{E} \int_0^t (u, h) ds. \quad (1.14)$$

If, in addition, $\mathbb{E}\|u_0\|_1^2 < \infty$, then

$$\mathbb{E}\|u(t)\|_1^2 + 2\nu \mathbb{E} \int_0^t |\Delta u(s)|_2^2 ds = \mathbb{E}\|u_0\|_1^2 + \mathfrak{B}_1 t + 2 \mathbb{E} \int_0^t (\nabla u, \nabla h) ds. \quad (1.15)$$

Time average. *There is $\gamma > 0$ depending only on $\{b_j\}$ such that*

$$\mathbb{P} \left\{ \sup_{t \geq 0} (\mathcal{E}_u(t) - (\mathfrak{B}_0 + 2\nu^{-1}|h|_2^2) t) \geq |u_0|_2^2 + \rho \right\} \leq e^{-\gamma\nu\rho}, \quad \rho > 0.$$

Exponential moment. *There is $c > 0$ not depending on ν , h and $\{b_j\}$ such that, if $\varkappa > 0$ and u_0 satisfy the inequalities*

$$\varkappa \sup_{j \geq 1} b_j^2 \leq c, \quad \mathbb{E} \exp(\varkappa\nu|u_0|_2^2) < \infty,$$

then, for some number $K = K(\nu, \varkappa, \mathfrak{B}_0, h)$, we have

$$\mathbb{E} \exp(\varkappa\nu|u(t)|_2^2) \leq e^{-\varkappa\nu^2 t} \mathbb{E} \exp(\varkappa\nu|u_0|_2^2) + K, \quad t \geq 0. \quad (1.16)$$

Higher Sobolev norms. *Suppose that $h \in H^k$ and $\mathfrak{B}_k < \infty$ for some integer $k \geq 1$. Then, for any $m \geq 1$ and $T \geq 1$, there is $C(k, m, T) > 0$ such that*

$$\mathbb{E} \sup_{0 \leq t \leq T} \mathcal{E}_u(k, t)^m \leq C(k, m, T) (1 + \nu^{-m(7k+2)} (\mathbb{E}|u_0|_2^{4m(k+1)} + 1)). \quad (1.17)$$

A straightforward consequence of the energy balance (1.14) and Gronwall's inequality is the exponentially fast stabilisation of the L^2 norms of solutions:

$$\mathbb{E}|u(t)|_2^2 \leq e^{-\nu t} \mathbb{E}|u_0|_2^2 + \nu^{-1} \mathfrak{B}_0 + \nu^{-2} |h|_2^2, \quad t \geq 0.$$

Combining this with (1.16) and (1.17), we conclude that, if all moments of $|u_0|_2^2$ are finite, then

$$\mathbb{E} \sup_{s \leq t \leq s+T} \mathcal{E}_u(k, t)^m \leq C'(k, m, T) \quad \text{for all } s \geq 0, m \in \mathbb{N}. \quad (1.18)$$

A proof of all these results, as well as of their counterparts for random kick-forces, can be found in Chapter 2 of [KS12].

The Markov process defined by solutions of Eqs. (0.1), (1.5) is time-homogeneous: the law at time t_2 of a solution $u(t)$ which takes a prescribed deterministic value v at $t = t_1 < t_2$ depends only on v and $t_2 - t_1$. Solutions of the kick-forced equation (0.1), (1.6) define an inhomogeneous Markov process, but its restriction to integer values of time $t \in \mathbb{Z}$ is a homogeneous Markov chain, and when studying Eqs. (0.1), (1.6) we usually restrict ourselves to integer t 's, see in [KS12]. Finally, solutions $u(t)$ of Eqs. (0.1), (1.7) do not define a Markov process, but their restrictions to $t \in \mathbb{Z}$ form a homogeneous Markov chain, which is the subject of our study when dealing with that equation, see [Shi15, KNS18].

2 Mixing

In the previous section, we discussed the existence, uniqueness, and regularity of the flow for the Navier–Stokes system subject to an external random force. Our next goal is to study its large-time asymptotics. As before, we shall mostly concentrate on the case of spatially regular white noise.

2.1 Existence of a stationary measure

We consider the Navier–Stokes system (0.1) with the right-hand side of the form (1.2), where $h \in H$ is a deterministic function and η has the form (1.5). Since $\mathfrak{B}_1 < \infty$, then ζ is a continuous functions of time with range in V .

Definition 2.1. A measure $\mu \in \mathcal{P}(H)$ is said to be *stationary* for the Navier–Stokes system if $\mathfrak{P}_t^* \mu = \mu$ for all $t \geq 0$.

By (1.12), (1.13), a measure $\mu \in \mathcal{P}(H)$ is stationary if and only if there is an H -valued random variable u_0 independent from ζ such that $\mathcal{D}(u_0) = \mu$, and one of the following equivalent properties is satisfied for the corresponding solution $u(t; u_0)$:

- (a) for all $g \in C_b(H)$ and $t \geq 0$, we have $\mathbb{E} g(u(t)) = \mathbb{E} g(u_0)$;
- (b) the measure $\mathcal{D}(u(t))$ coincides with μ for all $t \geq 0$.

A solution $u(t)$ of (0.1) as in (a) (or (b)) is called a *stationary solution*.

The existence of a stationary measure for the 2d Navier–Stokes system can be established with the help of the Borolyugov–Krylov argument, even though the first works dealing with both 2d and 3d cases used a different approach; see [VF88]. We refer the reader to the article [Fla94] a for self-contained and

simple proof of the existence of a stationary measure. In addition, Theorem 1.4 jointly with Fatou's lemma imply a priori estimates for any stationary measure of (0.1). A detailed proof of the following result can be found in [KS12, Chapter 2].

Theorem 2.2. *The Navier–Stokes system (0.1), (1.5) has at least one stationary measure μ_ν . It is supported by the space H^2 and satisfies the energy and enstrophy balance relations*

$$\nu \int_H |\nabla u|_2^2 \mu_\nu(du) = \frac{1}{2} \mathfrak{B}_0 + \int_H (h, u) \mu_\nu(du), \quad (2.1)$$

$$\nu \int_H |\Delta u|_2^2 \mu_\nu(du) = \frac{1}{2} \mathfrak{B}_1 + \int_H (\nabla h, \nabla u) \mu_\nu(du). \quad (2.2)$$

If, in addition, $h \in C^\infty$ and $\mathfrak{B}_k < \infty$ for all $k \geq 0$, then every stationary measure $\mu_\nu \in \mathcal{P}(H)$ is concentrated on infinitely smooth functions², and there are positive numbers \varkappa , C , and C_{km} not depending on ν such that, for any integers $m \geq 1$ and $k \geq 2$, we have

$$\int_H \exp(\varkappa \nu \|u\|_1^2) \mu_\nu(du) \leq C, \quad (2.3)$$

$$\nu^{m(7k+2)} \int_H \|u\|_k^{2m} \mu_\nu(du) \leq C_{km}. \quad (2.4)$$

Besides, for stationary solutions $u(t, x)$ of Eq. (0.1), estimate (1.17) implies uniform in $s \geq 0$ bounds of the form

$$\mathbb{E} \sup_{s \leq t \leq s+T} \|u(t)\|_k^{2m} \leq C(k, m, T) \nu^{-m(7k+2)};$$

see Corollary 2.4.13 in [KS12].

2.2 Uniqueness and exponential stability

In contrast to the existence of a stationary measure (which is established by rather soft tools), its uniqueness is a deep result that was proved thanks to the contribution of various research groups. It was first established in the case of spatially irregular white noise by Flandoli and Maslowski [FM95] and then extended to various types of regular noises in [KS00] and next in [EMS01, BKL02, KS01, Kuk02b, Mat02, KS02, KS03, Shi05, Oda08] (see Chapter 3 in [KS12] for more references). The following theorem summarises those results in the case of spatially regular white noise.

Theorem 2.3. *If the random process ζ in (1.5) satisfies*

$$b_j \neq 0 \quad \text{for all } j \in \mathbb{Z}_*^2, \quad (2.5)$$

then the problem (0.1), (1.5) has a unique stationary distribution $\mu_\nu \in \mathcal{P}(H)$. This measure possesses the following properties.

²This means that the μ_ν -measure of the space of C^∞ functions is equal to 1.

Exponential mixing. *There are positive numbers γ_ν and \varkappa_ν such that, for any locally Hölder-continuous function $g : V \rightarrow \mathbb{R}$ with at most exponential growth at infinity and any H -valued random variable u_0 independent from ζ , we have*

$$|\mathbb{E}g(u(t; u_0)) - \langle g, \mu_\nu \rangle| \leq C(\nu, g)e^{-\gamma_\nu t} \mathbb{E}e^{\varkappa_\nu |u_0|_2^2}, \quad t \geq 1, \quad (2.6)$$

where $C(\nu, g)$ depends on ν and a specific norm of g , but not on u_0 .

Convergence for observables. *If, in addition, $h \in H^k$ and $\mathfrak{B}_k < \infty$ for all k , then a similar convergence holds for any Hölder-continuous function g that is defined on a Sobolev space H^s of any finite order and has at most polynomial growth at infinity.*

Space homogeneity. *If, in addition to (2.5), $b_s \equiv b_{-s}$, then the measure μ_ν is space homogeneous (i.e. the space translations $H \ni u(x) \mapsto u(x + y)$, $y \in \mathbb{T}^2$, do not change it).*

The first assertion of the theorem implies that the Markov process in H which we discuss is *exponentially mixing in the Lipschitz-dual distance in the space H* , i.e. for each measure $\rho \in \mathcal{P}(H)$ with a finite second exponential moment we have

$$\|\mathfrak{P}_t^* \rho - \mu_\nu\|_{\text{Lip}(H)}^* \leq C_\nu e^{-\gamma'_\nu t} \quad \text{for } t \geq 0 \quad (2.7)$$

with some $C_\nu, \gamma'_\nu > 0$, where

$$\|\rho_1 - \rho_2\|_{\text{Lip}(H)}^* = \sup \langle \rho_1 - \rho_2, f \rangle, \quad \langle \rho, f \rangle = \int_H f(v) \rho(dv), \quad (2.8)$$

and the supremum is taken over all Lipschitz functions f on H whose norm and Lipschitz constant are bounded by one. If (2.7) holds, we also say that the stationary measure μ_ν and Eq. (0.1) are exponentially mixing (in H).

Inequality (2.6) expresses the property of convergence of the ensemble average of an observable g to its mean value with respect to the stationary measure. It applies to various physically relevant quantities, such as the energy $\frac{1}{2}|u|_2^2$, the enstrophy $\frac{1}{2}|\nabla u|_2^2 = \frac{1}{2}|\text{rot } u|_2^2$, and the correlation tensors $u_i(x)u_j(y)$, where $x, y \in \mathbb{T}^2$ are arbitrary points.

Similar results hold for solutions of Eqs. (0.1), (1.6) if in (2.6) we replace $t \geq 1$ by $t \in \mathbb{N}$.

In Theorem 2.3, the rate of convergence depends on the viscosity ν : when the latter decreases, the attractor of the unperturbed problem becomes larger and more chaotic, and it seems to be a very complicated task to establish a uniform convergence to the limiting measure. On the other hand, it is not difficult to prove that, when $\nu > 0$ is fixed, the non-degeneracy condition (2.5) can be relaxed, requiring only that the noise should act directly on the (finitely many³) determining modes of the dynamics. A challenging problem, important,

³The fact that dissipative parabolic-type PDEs have only finitely many determining modes goes back to the paper [FP67] and implies, in particular, that the global attractor is finite-dimensional.

in particular, for numerical simulations, is to prove that the mixing property remains true under a weaker hypothesis on the noise, allowing for the noise's localisation in a part of either the physical or the Fourier spaces, so that the determining modes of the dynamics do not necessarily belong to the region of the phase space affected by the noise. Propagation of the randomness then may take place due to the “mixing” properties of the (deterministic) Navier–Stokes flow. Uniqueness and mixing of the random flow in this situation are mostly established in the case when the deterministic component h of the random force is zero.

We begin with the case when the random force is localised in the Fourier space. In [HM06, HM11] Hairer and Mattingly obtained the following result:

Theorem 2.4. *Let the random force f have the form (1.2), (1.5) with $h = 0$ and with η satisfying*

$$b_j \neq 0 \quad \text{if and only if } j \in \mathcal{I}, \quad (2.9)$$

where \mathcal{I} is a finite subset of \mathbb{Z}_*^2 such that any vector in \mathbb{Z}^2 can be represented as an integer linear combination of the elements of \mathcal{I} , and \mathcal{I} contains at least two vectors of different length. Then Eq. (0.1) is exponentially mixing in H .

A key ingredient of the proof is an infinite-dimensional version of the Malliavin calculus, which uses the white noise structure of the random perturbation. Note that the result above does not imply the assertion of Theorem 2.3 since now the set of modes $j \in \mathbb{Z}_*^2$, excited by the random force, must be finite.

The recent paper [KNS18] deals with the case when the noise is a piecewise independent random process of the form (1.7). The main result of [KNS18] is an abstract theorem, establishing the mixing property for a large class of nonlinear PDE, perturbed by bounded random forces of the form (1.2), (1.7); its proof relies on the method of optimal control and a variant of the Nash–Moser scheme. In particular, the theorem in [KNS18] applies to the Navier–Stokes system, perturbed by a bounded red noise as in Example 1.3:

Theorem 2.5. *Let the random force f have the form (1.2), in which $h = 0$ and η is given (1.8), (1.9), with the set of coefficients b_j satisfying (2.9). Assume that $|c_n| \leq Cn^{-q}$ for all $n \geq 1$ with some $q > 1$, and that the law of the random variables ξ_l and ξ_l^n has the form $\lambda = p(r) dr$ with $p \in C_0^1(-1, 1)$, $p(0) \neq 0$. Then Eq. (0.1) is exponentially mixing in H .*

The situation in which the random perturbation is localised in the physical space was studied in [Shi15] (see also [Shi18] for the case of a boundary noise). To formulate the corresponding result, we fix a non-empty open set $Q \subset \mathbb{R} \times \mathbb{T}^2$ whose closure is contained in $(0, 1) \times \mathbb{T}^2$ and denote by $\{\varphi_j\} \subset H^1(Q, \mathbb{R}^2)$ an orthonormal basis in $L^2(Q, \mathbb{R}^2)$. Considering again the random force (1.2), we assume that $h \equiv 0$, and η is a piecewise independent random process of the form (1.7), with

$$\eta^k(t, x) = \sum_{j=1}^{\infty} b_j \xi_{jk} \psi_j(t, x).$$

Here $\psi_j = \chi\varphi_j$, where $\chi \in C_0^\infty(Q)$ is a nonzero function, $\{b_j\}$ are real numbers such that $\sum_j b_j \|\psi_j\|_1 < \infty$, and ξ_{jk} are independent random variables whose laws have the form $\lambda_j = p_j(r) dr$, where $p_j \in C_0^1(-1, 1)$ and $p_j(0) \neq 0$. Thus, the random force entering the right-hand side of (0.1) is bounded and space-time localised in Q . The following theorem is the main result of [Shi15].

Theorem 2.6. *In addition to the above hypotheses, assume that $b_j \neq 0$ for all $j \geq 1$. Then, for any $\nu > 0$, Eq. (0.1) has a unique stationary measure, which is exponentially mixing in H .*

3 Consequences of mixing

The results of the previous section concern the evolution of the mean values of observables and the laws of solutions under the stochastic Navier–Stokes flow. This section is devoted to studying the typical behaviour of individual trajectories. We will assume that the random force in equation (0.1) is such that the equation is exponentially mixing, either for $t > 0$, or for $t \in \mathbb{N}$. That is, either the assumptions of Theorems 2.3, 2.4, 2.5 or 2.6 hold, or f is a kick-force of the form (1.2), (1.6), where (2.5) holds. We state the results for the case when Theorem 2.3 applies. Situation with the other cases is very similar; e.g., see [KS12] for the case of the kick-forced equations.

3.1 Ergodic theorems

We consider the Navier–Stokes system (0.1), (1.2) in which $h \in H$ is a fixed function and η is given by (1.5). Let us denote⁴ by \mathcal{H} the space of locally Hölder-continuous functions $g : V \rightarrow \mathbb{R}$ with at most exponential growth at infinity. The following result established in [Kuk02a, Shi06] (see also Section 4.1.1 in [KS12]) shows that the time average of a large class of observables converges to their mean value with respect to the stationary measure.

Theorem 3.1 (Strong law of large numbers). *Under the hypotheses of Theorem 2.3, for any $\gamma \in [0, 1/2)$, any non-random $u_0 \in H$ and $g \in \mathcal{H}$, with probability 1 we have*

$$\lim_{t \rightarrow \infty} t^\gamma \left(\frac{1}{t} \int_0^t g(u(s; u_0)) ds - \langle g, \mu_\nu \rangle \right) = 0. \quad (3.1)$$

Convergence (3.1) remains true for random initial functions that are independent from η and have a finite exponential moment. Moreover, some further analysis shows that the law of iterated logarithm (LIL) is also valid. In particular, the number γ in (3.1) characterizing the rate of convergence to the mean value cannot be taken to be equal to $1/2$; see Section 4.1.2 in [KS12].

⁴To avoid unimportant complications, we do not give an exact definition of the space \mathcal{H} , referring the reader to Section 4.1 in [KS12].

We now turn to the central limit theorem (CLT). To this end, given an observable $g \in \mathcal{H}$ with zero mean value with respect to μ_ν , we denote

$$\sigma_g^2 = 2 \int_H \int_0^\infty (\mathfrak{P}_t g)(v) dt g(v) \mu_\nu(dv).$$

In view of inequality (2.6) and the assumption that $\langle g, \mu_\nu \rangle = 0$, the function $\mathfrak{P}_t g$ decays exponentially to zero as $t \rightarrow \infty$, and it is not difficult to prove that, under the hypotheses of Theorem 2.3, the number σ_g^2 is well defined and positive for any non-constant g ; see Proposition 4.1.4 in [KS12]. A proof of the following result can be found in [Kuk02a, Shi06] (see also Section 4.1.3 in [KS12]).

Theorem 3.2 (Central limit theorem). *Under the hypotheses of Theorem 2.3, for any non-random $u_0 \in H$ and any non-constant $g \in \mathcal{H}$ satisfying $\langle g, \mu_\nu \rangle = 0$, we have*

$$\mathcal{D} \left(\frac{1}{\sqrt{t}} \int_0^t g(u(s; u_0)) ds \right) \rightarrow \mathcal{N}_{\sigma_g} \quad \text{as } t \rightarrow \infty, \quad (3.2)$$

where \mathcal{N}_σ stands for the centred normal law on \mathbb{R} with a variance $\sigma^2 > 0$, and \rightarrow stands for the weak convergence of measures.

Let us emphasise that σ_g is expressed in terms of the stationary measure μ_ν and does not depend on u_0 . Furthermore, convergence (3.2) is equivalent to the relation

$$\lim_{t \rightarrow \infty} \mathbb{P} \left\{ \frac{1}{\sqrt{t}} \int_0^t g(u(s)) ds \in \Gamma \right\} = \mathcal{N}_{\sigma_g}(\Gamma), \quad (3.3)$$

where $\Gamma \subset \mathbb{R}$ is an arbitrary Borel set whose boundary has zero Lebesgue measure.

3.2 Random attractors

Another important object characterising the large-time behaviour of trajectories, is the random attractor. There are many definitions of this object, and in the context of random dynamical systems most of them deal with the concept of pullback attraction. The meaning of the latter is that, for some fixed observation time, the distance between trajectories of the system and the attractor decreases when the moment of beginning of the observation goes to $-\infty$. In this section, we discuss a concept of attractor that possesses an attraction property forward in time and is closely related to the unique stationary measure constructed in Theorem 2.3. To simplify notation, we fix the viscosity $\nu > 0$ and do not follow the dependence of various objects on it. Moreover, we assume that the Brownian motions $\{\beta_j\}$ entering (1.5) are two-sided⁵ and denote by $(\Omega, \mathcal{F}, \mathcal{F}_t, \mathbb{P})$ the canonical filtered probability space associated with the process ζ ; see Section 2.4 in [KS91]. In particular, the measurable space (Ω, \mathcal{F}) may be chosen to coincide with the Fréchet space $C(\mathbb{R}, V)$ of continuous functions from \mathbb{R} to V , with the

⁵That is, they are defined for all $t \in \mathbb{R}$ and vanish at $t = 0$.

topology of uniform convergence on bounded intervals and the corresponding Borel sigma-algebra.

Let us denote by $\varphi_t(\omega) : H \rightarrow H$ the random flow generated by the Navier–Stokes system with spatially regular white noise (1.5). Thus, for any initial function $u_0 \in H$, the solution $u(t; u_0)$ of (0.1), (0.2), (1.5) is given by $\varphi_t(\omega)u_0$. The family $\{\varphi_t(\omega), \omega \in \Omega, t \geq 0\}$ possesses the perfect co-cycle property. Namely, let $\theta_t : \Omega \rightarrow \Omega$ be the shift operator taking $\omega(\cdot)$ to $\omega(\cdot + t)$. Then, there is a set of full measure $\Omega_0 \in \mathcal{F}$ such that, for any $\omega \in \Omega_0$, we have

$$\varphi_{t+s}(\omega)u_0 = \varphi_t(\theta_s(\omega))\varphi_s(\omega)u_0 \quad \text{for all } t, s \geq 0, u_0 \in H.$$

The following result is established in [Led86, Le 87, Cra91] (see also Theorem 4.2.9 in [KS12]) in the context of general random dynamical systems.

Theorem 3.3. *Under the hypotheses of Theorem 2.3, for any sequence $\{t_k\}$ going to $+\infty$ the limit*

$$\mu_\omega = \lim_{k \rightarrow \infty} \varphi_{t_k}(\theta_{-t_k}(\omega))_* \mu \tag{3.4}$$

exists in the weak topology of $\mathcal{P}(\Omega)$ for almost every $\omega \in \Omega$. Moreover, the following properties hold:

Uniqueness. *If $\{t'_k\}$ is another sequence going to $+\infty$ and $\{\mu'_\omega\}$ is the corresponding limit, then $\mu_\omega = \mu'_\omega$ almost everywhere.*

Reconstruction. *For any $\Gamma \in \mathcal{B}(H)$, the mapping $\omega \mapsto \mu_\omega(\Gamma)$ is measurable, and μ can be reconstructed by the formula $\mu(\Gamma) = \mathbb{E}\mu_\omega(\Gamma)$.*

We now describe a random attractor associated with μ . To this end, we fix any sequence $\{t_k\}$ going to $+\infty$ and denote by Ω_0 the set of full measure on which limit (3.4) exists. The almost sure existence of the limit in (3.4) implies that Ω_0 is \mathbb{P} -invariant under θ_t ; that is, $\mathbb{P}(\Omega_0 \Delta (\theta_t(\Omega_0))) = 0$ for any $t \in \mathbb{R}$, where Δ stands for the symmetric difference of two sets. We define μ_ω by (3.4) on the set of full measure Ω_0 and denote by \mathcal{A}_ω the support of μ_ω for $\omega \in \Omega_0$, while we set $\mathcal{A}_\omega = \emptyset$ on the complement of Ω_0 . The measurability property of μ mentioned in Theorem 3.3 implies that $\{\mathcal{A}_\omega, \omega \in \Omega\}$ is also measurable in the sense that, for any $u \in H$, the function $\omega \mapsto d_H(u, \mathcal{A}_\omega)$ is measurable. The following result is established in [KS04] (see also Section 4.2 in [KS12]).

Theorem 3.4. *Under the hypotheses of Theorem 2.3, the following properties hold:*

Invariance. *For any $t \geq 0$ and almost every $\omega \in \Omega$, we have $\varphi_t(\omega)\mathcal{A}_\omega = \mathcal{A}_{\theta_t(\omega)}$.*

Attraction. *For any $u \in H$, the functions $\omega \mapsto d_H(\varphi_t(\omega)u, \mathcal{A}_{\theta_t(\omega)})$ converge to zero in probability as $t \rightarrow \infty$. That is, for any $\varepsilon > 0$, we have*

$$\mathbb{P}\{d_H(\varphi_t(\omega)u, \mathcal{A}_{\theta_t(\omega)}) \geq \varepsilon\} \rightarrow 0 \quad \text{as } t \rightarrow \infty.$$

Minimality. *If $\{\mathcal{A}'_\omega, \omega \in \Omega\}$ is another measurable family of closed subsets that satisfies the first two properties, then $\mathcal{A}_\omega \subset \mathcal{A}'_\omega$ for almost every $\omega \in \Omega$.*

3.3 Dependence on parameters and stability

We now investigate how the laws of trajectories of (0.1) vary with parameters. To simplify the presentation, we shall consider only the dependence on the random forcing, assuming that it is a spatially regular white noise. However, similar methods can be used to study the dependence on other parameters, such as the viscosity, as well as the relationship between stationary measures corresponding to various types of random forcings.

We thus assume that X is a metric space and that the right-hand side in (0.1) has the form (1.2), (1.5), where $h \in H$ is fixed and b_j 's are continuous functions of $a \in X$, satisfying

$$\sup_{a \in X} \sum_{j \in \mathbb{Z}_*^2} b_j(a)^2 |j|^2 < \infty.$$

A proof of the following result can be found in Section 4.3.1 of [KS12].

Theorem 3.5. *In addition to the above hypotheses, let*

$$b_j(\hat{a}) \neq 0 \quad \text{for all } j \in \mathbb{Z}_*^2, \quad (3.5)$$

where $\hat{a} \in X$ is a fixed point, and let $\{\mu^a\}$ be a family of stationary measures⁶ of Eq. (0.1) with the right-hand side corresponding to the value $a \in X$ of the parameter. Then $\mu^a \rightarrow \mu^{\hat{a}}$ as $a \rightarrow \hat{a}$, and for any compact subset $\Lambda \subset \mathcal{P}(H)$ there exists a continuous function $A_\Lambda(\rho) > 0$ going to zero with ρ such that

$$\sup_{t \geq 0} \|\mathfrak{P}_t^*(a)\lambda_1 - \mathfrak{P}_t^*(\hat{a})\lambda_2\|_L^* \leq A_\Lambda(\|\lambda_1 - \lambda_2\|_L^* + d_X(a, \hat{a})) \quad \text{for } \lambda_1, \lambda_2 \in \Lambda, a \in X,$$

where $\mathfrak{P}_t^*(a)$ denotes the Markov semigroup for (0.1) corresponding to the parameter a .

We emphasise that by this result the law of a solution for (0.1) depends on the parameters of the random force f continuously and uniformly in time; cf. Theorem 6.1 below.

4 Large deviations

Having discussed the typical behaviour of trajectories, we now turn to a description of probabilities of rare events. Two different asymptotics will be studied: deviations of the time-average of observables from their ensemble average as $t \rightarrow \infty$, and deviations from the limiting dynamics in the small noise regime. In the latter setting, we shall only discuss the behaviour of a stationary distribution, since in the case of an additive noise the asymptotics of trajectories with given initial data is a simple result that follows immediately from the large deviations principle (LDP) for Gaussian random variables. We refer the reader to the paper [CM10] and the references therein for this type of results for stochastic PDEs with multiplicative noise.

⁶Existence of a measure μ^a is guaranteed by Theorem 2.2, and Theorem 2.3 implies that the measure $\mu^{\hat{a}}$ is uniquely determined.

4.1 Donsker–Varadhan type large deviations

We first describe the general idea. To this end, let us note that, for an arbitrary function $g \in \mathcal{H}$ (whose mean value is not necessarily zero), we can rewrite the convergence (3.3) in the form

$$\lim_{t \rightarrow \infty} \mathbb{P} \left\{ \frac{1}{t} \int_0^t g(u(s)) \, ds \in \langle g, \mu_\nu \rangle + \frac{\Gamma}{\sqrt{t}} \right\} = \mathcal{N}_{\sigma_g}(\Gamma), \quad (4.1)$$

Thus, the CLT can be interpreted as a description of the probabilities of small deviations of the time average of observables from their mean value. The goal of the theory of large deviations is to describe the probabilities of order 1 deviations (when Γ/\sqrt{t} in (4.1) is replaced by Γ). This type of results were first obtained by Donsker and Varadhan in the case of finite-dimensional diffusion processes [DV75] and later extended to many other situations. In the context of the Navier–Stokes equations, the theory was developed in the case of random kicks [JNPS15, JNPS17] (see also the recent papers [MN17, Ner17] devoted to spatially regular white noise), and we now describe the main achievements.

Let us consider the Navier–Stokes system (0.1), (1.2), where $h \in H$ is a deterministic function, and the random forcing η is given by (1.6). In this case, the trajectories of (0.1) have jumps at integer times, and we normalise them to be right-continuous in time. Setting $u_k = u(k)$, we see that the random sequence $\{u_k\}$ satisfies the relation

$$u_k = S(u_{k-1}) + \eta^k, \quad k \geq 1, \quad (4.2)$$

where $S : H \rightarrow H$ denotes the time-1 shift along trajectories of Eq. (0.1) with $f = h$. Equation (4.2) defines a discrete-time Markov process in H , and we use the notation introduced in Section 1.2 to denote the related objects, replacing t with k . We thus write $P_k(v, \Gamma)$, \mathfrak{P}_k , and \mathfrak{P}_k^* .

To formulate the result on LDP, we shall need some hypotheses on the kicks η^k . Namely, we shall assume that they satisfy the following condition, in which $\{e_j\}$ is the trigonometric basis in H defined by (1.3).

Structure of the noise. *The random kicks η^k have the form*

$$\eta^k(x) = \sum_{j \in \mathbb{Z}_*^2} b_j \xi_{jk} e_j(x),$$

where b_j are some non-zero numbers satisfying $\mathfrak{B}_1 < \infty$ and $\{\xi_{jk}\}$ are independent random variables whose laws possess C^1 -smooth positive densities ρ_j with respect to the Lebesgue such that $\int_{\mathbb{R}} |\rho'_j(r)| \, dr \leq 1$ for any j .

This condition ensures that the Markov process $\{u_k\}$ associated with (0.1) has a unique stationary measure μ_ν for any $\nu > 0$, and the CLT holds for any Holder-continuous functional $f : H \rightarrow \mathbb{R}$ with at most exponential growth at infinity (cf. Theorem 2.3 and 3.2).

We now introduce the *occupation measures*

$$\mu_n^\omega = \frac{1}{n} \sum_{k=0}^{n-1} \delta_{u_k},$$

where $\delta_v \in \mathcal{P}(H)$ stands for the Dirac mass at $v \in H$, and $\{u_k\}$ is a trajectory of (4.2). Let us recall that the space $\mathcal{P}(H)$ is endowed with the topology of weak convergence. We shall say that a mapping $I : \mathcal{P}(H) \rightarrow [0, +\infty]$ is a *good rate function* if its sub-level set $\{\lambda \in \mathcal{P}(H) : I(\lambda) \leq c\}$ is compact for any $c \geq 0$. The following result is established in Section 2.2 of [JNPS17] (see also [JNPS15] for the case of bounded kicks).

Theorem 4.1. *Suppose that the above hypothesis on the structure of the noise is satisfied. Then the following assertions hold for any $\nu > 0$:*

Pressure. *For any $g \in C_b(H)$ and any deterministic initial function $u_0 \in H$, there is a finite limit*

$$Q(g) = \lim_{n \rightarrow \infty} n^{-1} \log \mathbb{E} \exp \left\{ \sum_{k=0}^{n-1} g(u_k) \right\}$$

that is independent from u_0 . Moreover, Q is a 1-Lipschitz function satisfying the relation $Q(g + C) = Q(g) + C$ for any $C \in \mathbb{R}$.

Rate function. *The Legendre transform $I : \mathcal{P}(H) \rightarrow [0, +\infty]$ of Q defined by*

$$I(\lambda) = \sup_{g \in C_b(H)} (\langle g, \lambda \rangle - Q(g))$$

is a convex good rate function that can be represented by the Donsker-Varadhan relation

$$I(\lambda) = \sup_{g \geq 1} \int_H \log \frac{g}{\mathfrak{F}_1 g} d\lambda,$$

where the supremum is taken over all functions $g \in C_b(H)$ minorised by 1.

LDP. *For any random initial function u_0 independent from $\{\eta^k\}$ such that $\mathbb{E} \exp(\delta |u_0|_2^2) < \infty$ for some $\delta > 0$, and any Borel set $\Gamma \subset \mathcal{P}(H)$, we have*

$$-I(\dot{\Gamma}) \leq \liminf_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\{\mu_n \in \Gamma\} \leq \limsup_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\{\mu_n \in \Gamma\} \leq -I(\bar{\Gamma}),$$

where $\dot{\Gamma}$ and $\bar{\Gamma}$ denote the interior and closure of Γ , respectively, and $I(A)$ is the infimum of I on the set A .

Assuming faster decay for the coefficients b_j and considering the Navier-Stokes system in higher Sobolev spaces, it is not hard to show that a similar result holds in the case when H is replaced by H^s . Next, application of the standard techniques of the theory of large deviations implies that, for any Hölder-continuous function $g : H^s \rightarrow \mathbb{R}$ with moderate growth at infinity, the time average $n^{-1} \sum_{k=0}^{n-1} g(u_k)$ satisfies the LDP with a good rate function $I_g : \mathbb{R} \rightarrow \mathbb{R}$ that can be expressed in terms of I by the relation (cf. Section 1.3 in [JNPS15])

$$I_g(r) = \inf \{I(\lambda) : \lambda \in \mathcal{P}(H), \langle g, \lambda \rangle = r\}.$$

4.2 Vanishing noise limit

We now go back to the Navier–Stokes system (0.1) with a spatially regular white noise and discuss the behaviour of the unique stationary measure as the stochastic component of the noise goes to zero. Namely, let us assume that the external force in (0.1) has the form

$$f(t, x) = h(x) + \sqrt{\varepsilon} \eta(t, x), \quad (4.3)$$

where $h \in H$ is a fixed function, $\varepsilon > 0$ is a small parameter, and η is given by (1.5). We assume that the coefficients b_j are non-zero (and satisfy the inequality $\mathfrak{B}_1 < \infty$), so that for any $\nu > 0$ and $\varepsilon > 0$ there is a unique stationary measure $\mu_\nu^\varepsilon \in \mathcal{P}(H)$. We are interested in the behaviour of μ_ν^ε for a fixed ν , so we drop ν from the notation in this section and write simply μ^ε . In what follows, we assume that the following hypothesis is satisfied.

Global asymptotic stability. *The flow of the unperturbed Navier–Stokes system (0.1), corresponding to $f = h$, has a unique fixed point $\hat{u} \in H$, which is globally asymptotically stable in the sense that any other trajectory converges to it as $t \rightarrow +\infty$.*

Notice that this condition is satisfied with $\hat{u} = 0$ if the deterministic part of the force (1.2) vanishes. A simple argument based on the uniqueness of a stationary distribution for the limiting equation implies that $\{\mu^\varepsilon\}$ converges weakly to the Dirac mass concentrated at \hat{u} . In fact, much more detailed information about that convergence is available. Namely, let us denote by $S_t(u_0, f)$ the solution of (0.1), (0.2) and introduce a *quasi-potential* by the relation

$$\mathcal{V}(v) = \liminf_{r \rightarrow 0} \left\{ \frac{1}{2} \int_0^T \|f(s)\|_b^2 ds : T > 0, f \in L^2(J_T, H), |S_T(\hat{u}, f) - v|_2 \leq r \right\},$$

where the infimum is taken over all T and f for which the inequality holds, and we set $\|g\|_b^2 = \sum_j b_j^{-2}(g, e_j)^2$. A proof of the following result can be found in [Mar17a] (see also [BC17] for the case of spatially irregular noise).

Theorem 4.2. *Suppose that global asymptotic stability for the limiting dynamics holds, and the coefficients $\{b_j\}$ entering (1.5) are non-zero. Then the function $\mathcal{V} : H \rightarrow [0, +\infty]$ has compact sub-level sets in H , vanishes only at the point \hat{u} , and controls the LDP for the family $\{\mu^\varepsilon\}$; that is, for any $\Gamma \in \mathcal{B}(H)$, we have*

$$-\inf_{v \in \Gamma} \mathcal{V}(v) \leq \liminf_{\varepsilon \rightarrow 0^+} \varepsilon \log \mu^\varepsilon(\Gamma) \leq \limsup_{\varepsilon \rightarrow 0^+} \varepsilon \log \mu^\varepsilon(\Gamma) \leq -\inf_{v \in \Gamma} \mathcal{V}(v).$$

The facts that \mathcal{V} vanishes only at \hat{u} and has compact sub-level sets imply that $\mu^\varepsilon(H \setminus B) \sim e^{-c(B)/\varepsilon}$ as $\varepsilon \rightarrow 0$, where $B \subset H$ is an arbitrary ball around \hat{u} , and $c(B) > 0$ is a number. Hence, the LDP gives an estimate for the rate of concentration of μ^ε around \hat{u} . Let us also note that, in the case when the limiting dynamics is not globally asymptotically stable, it is still possible to

prove that the family $\{\mu^\varepsilon\}$ is exponentially tight, but the validity of the LDP is not known to hold (see Open problem 3 in Section 7). However, more detailed information on the limiting dynamics would be sufficient to get the LDP. For instance, this is the case when there are finitely many stationary points, and the unperturbed dynamics possesses a global Lyapunov function. We refer the reader to [Mar17b] for detail.

5 Inviscid limit

If the random force f in Eq. (0.1) has the form (4.3) with $h = 0$ and with η as in (1.5), where $\nu > 0$ is fixed and $\varepsilon \rightarrow 0$, then the corresponding stationary measure μ_ν^ε converges to the delta-measure at $0 \in H$, and the results of Section 4.2 describe this convergence in more detail. Now assume that ε and ν both go to zero in such a way that $\varepsilon = \nu^a$ with some $a > 0$. Applying to $\mu_\nu^{\nu^a}$ relation (2.1) with $h = 0$ and $b_j := \sqrt{\varepsilon} b_j$, $j \in \mathbb{Z}_*^2$, we get

$$\int_H |\nabla u|_2^2 \mu_\nu^{\nu^a}(du) = \frac{\mathfrak{B}_0}{2} \nu^{a-1}.$$

From this we conclude that the measure $\mu_\nu^{\nu^a}$ may have a non-trivial limit as $\nu \rightarrow 0$ only if $a = 1$. Then Eq. (0.1) becomes

$$\partial_t u + \langle u, \nabla \rangle u - \nu \Delta u + \nabla p = \sqrt{\nu} \eta(t, x), \quad \operatorname{div} u = 0. \quad (5.1)$$

Assume that (2.5) holds. Then a stationary measure μ_ν^ν of (5.1) is unique. For short, we re-denote it as μ^ν .⁷

5.1 Properties of μ^ν , independent from ν .

In this subsection, we discuss a number of properties of the stationary measures μ^ν which hold uniformly in ν . They depend only on the quantities \mathfrak{B}_0 and \mathfrak{B}_1 and indicate certain universal properties of the statistical equilibria of Eq. (5.1).

Relations (2.1) and (2.2) with $h = 0$ and $\mathfrak{B}_0 =: \nu \mathfrak{B}_0$, $\mathfrak{B}_1 =: \nu \mathfrak{B}_1$ imply that, uniformly in ν ,

$$\mathbb{E}^{\mu^\nu} |\nabla u|_2^2 = \frac{1}{2} \mathfrak{B}_0, \quad \mathbb{E}^{\mu^\nu} |\Delta u|_2^2 = \frac{1}{2} \mathfrak{B}_1. \quad (5.2)$$

Since $|\nabla u|_2^2 \leq |u|_2 |\Delta u|_2$, then $\mathbb{E}^{\mu^\nu} |\nabla u|_2^2 \leq (\mathbb{E}^{\mu^\nu} |u|_2^2)^{1/2} (\mathbb{E}^{\mu^\nu} |\Delta u|_2^2)^{1/2}$. It follows that $\mathfrak{B}_0^2 / 2 \mathfrak{B}_1 \leq \mathbb{E} |u^\nu|_2^2 \leq \frac{1}{2} \mathfrak{B}_0$, so that the averaged kinetic energy is bounded below and above. Moreover, uniformly in ν , the measures μ^ν satisfy (2.3) with $\nu := 1$ (this follows immediately from the proof in [KS12]). Consider a stationary solution $u^\nu(t, x)$, corresponding to μ^ν . Then $\mathbb{E} |u^\nu|_2^2 = \mathbb{E}^{\mu^\nu} |u|_2^2$, and the

⁷Since now f has the form (1.2) with $h = 0$, then the uniqueness of the stationary measure also follows from (2.9), but we need (2.5) for the validity of some results in this section.

Reynolds number of u^ν is

$$\text{Re}_\nu = \frac{[u^\nu][x]}{\nu} = \frac{(\mathbb{E}|u^\nu|_2^2)^{1/2} \cdot 1}{\nu} \sim \nu^{-1}$$

($[\cdot]$ stands for the characteristic size of a variable), while the averaged kinetic energy $\frac{1}{2}(\mathbb{E}|u^\nu|_2^2)$ is of order one. So when $\nu \rightarrow 0$ the solutions u^ν describe space-periodic stationary 2d turbulence.

Assume that $b_s \equiv b_{-s}$, so the measures μ^ν are space-homogeneous, and that b_s decay sufficiently fast when $|s| \rightarrow \infty$. In this case, as it is shown in [KP05, KS12], the measures μ^ν possess additional properties. Namely, let $g(r)$ be any continuous function, having at most a polynomial growth at infinity. Then, denoting $v = \text{rot } u$, we have the following *balance relation*, valid for all $\nu > 0$:

$$\mathbb{E}^{\mu^\nu} g(v(t, x)) |\nabla v(t, x)|^2 = \frac{1}{2} (2\pi)^{-2} \mathfrak{B}_1 \mathbb{E}^{\mu^\nu} g(v(t, x)) \quad (5.3)$$

(by the homogeneity the relation does not depend on x). Since $|\Delta u|_2^2 = |\nabla v|_2^2$, then relations (5.2) and the translational invariance of μ^ν imply that

$$\mathbb{E}^\mu |\nabla v(t, x)|^2 = \frac{1}{2} (2\pi)^{-2} \mathfrak{B}_1.$$

So (5.3) means that the random variables $|\nabla v(t, x)|^2$ and $g(v(t, x))$ are uncorrelated, for any continuous function g as above and any (t, x) .

The balance relations (5.3) admit a surprising reformulation. For any $\tau \in \mathbb{R}$ denote by $\Gamma_\tau(\omega)$ the random curve $\{x \in \mathbb{T}^2 : v^\omega(t, x) = \tau\}$ (it is well defined for a.a. τ and ω , if b_s decays fast enough). Then

$$\mathbb{E}^{\mu^\nu} \int_{\Gamma_\tau(\omega)} |\nabla v^\omega| \, d\ell = \frac{1}{2} (2\pi)^{-1} \mathfrak{B}_1 \mathbb{E}^{\mu^\nu} \int_{\Gamma_\tau(\omega)} |\nabla v^\omega|^{-1} \, d\ell, \quad \text{for a.a. } \tau, \quad (5.4)$$

where $d\ell$ is the length element on $\Gamma_\tau(\omega)$, and the existence of the integrals in the l.h.s. and the r.h.s. for a.a. τ is a part of the assertion. This is the *co-area form of the balance relations*. Besides, relations (5.3) imply the following point-wise exponential estimates

$$\mathbb{E}^{\mu^\nu} \left(e^{\sigma|v(t,x)|} + e^{\sigma|u(t,x)|} + e^{\sigma|\nabla u(t,x)|^{1/2}} \right) \leq K \quad \forall x, \quad (5.5)$$

valid uniformly in ν , where the positive constants σ and K depend only on the first few numbers \mathfrak{B}_j ; see [KS12].

5.2 Inviscid limit

Since estimates (5.2) hold uniformly in ν , the family of measures $\{\mu_\nu, 0 < \nu \leq 1\}$ is tight in $H^{2-\epsilon}$ and, by Prokhorov's theorem, relatively compact in the space $\mathcal{P}(H^{2-\epsilon})$, for any $\epsilon > 0$. So any sequence of measures $\{\mu^{\nu'_j}, \nu'_j \rightarrow 0\}$, contains a weakly converging subsequence:

$$\mu^{\nu'_j} \rightarrow \mu \quad \text{weakly in } \mathcal{P}(H^{2-\epsilon}). \quad (5.6)$$

Relations (5.2) immediately imply that

$$\mathbb{E}^\mu |\nabla u|_2^2 = \frac{1}{2} \mathfrak{B}_0, \quad \mathbb{E}^\mu |\Delta u|_2^2 \leq \frac{1}{2} \mathfrak{B}_1, \quad \mathfrak{B}_0^2 / 2\mathfrak{B}_1 \leq \mathbb{E}^\mu |u|^2 \leq \frac{1}{2} \mathfrak{B}_0, \quad (5.7)$$

so μ is supported by the space H^2 . More delicate analysis of the convergence (5.6) shows that

$$\mu(K) = 1, \quad \text{where } K = \{u \in H^1 : \text{rot } u \in L_\infty\}, \quad (5.8)$$

see [GSV15]. Moreover, the measure μ is invariant for the deterministic equation (5.1)| $_{\nu=0}$, i.e. for the free 2d Euler equation

$$\partial_t u + (u \cdot \nabla)u + \nabla p = 0, \quad \text{div } u = 0. \quad (5.9)$$

See [Kuk04, KS12].⁸ The limit (5.6) is the *inviscid limit* for the (properly scaled) stochastic 2d NSE. In view of what has been said at the beginning of the previous subsection, the inviscid limit measures μ describe the statistic of space-periodic stationary 2d turbulence. We summarise the results concerning this limit in a theorem:

Theorem 5.1. (1) *Any inviscid limit measure μ satisfies (5.7), (5.8) and is invariant for Eq. (5.9).*

(2) *If $b_s \equiv b_{-s}$ and $|b_s|$ decay sufficiently fast as $|s| \rightarrow \infty$, then the measure μ is space-homogeneous and satisfies (5.5), where \mathbb{E}^{μ^ν} should be replaced by \mathbb{E}^μ .*

The last assertion follows from (5.5), convergence (5.6) and Fatou's lemma. Convergence (5.6) does not allow to pass to the limit in (5.3) and (5.4), and we do not know if the balance relations hold for the inviscid limit measures μ . Relations (5.7) imply that the limiting measures μ are non-trivial in the sense that they do not equal the delta-measure in the origin. In fact, they are non-degenerate in a much stronger sense. To state the corresponding result, we denote

$$E(u) = \frac{1}{2} |u|_2^2, \quad E_1(u) = \frac{1}{2} |\text{rot } u|_2^2 = \frac{1}{2} |\nabla u|_2^2,$$

and call a real analytic function $f(r)$ *admissible* if $f''(r)$ has at most a polynomial growth as $|r| \rightarrow \infty$ and is bounded from below (e.g., f is a polynomial of the form $f = r^{2m} + \dots$, or any trigonometric polynomial).

Theorem 5.2 ([Kuk08, KS12]). (1) *The inviscid limits μ are such that the push-forward measures $E_*\mu$ and $(E_1)_*\mu$ are absolutely continuous with respect to the Lebesgue measure on \mathbb{R} .*

⁸The measure μ is supported by the space K , on which the flow of the Euler equation is well defined in the sense of Yudovich and is continuous in the weak* topology of that set, see in [GSV15]. Before (5.8) was obtained, an additional construction, suggested in [Kuk04], was used to explain in which sense the limiting measure μ is invariant for (5.9).

(2) If $\mathfrak{B}_2 < \infty$ and $b_s \equiv b_{-s}$, then for any $d \in \mathbb{N}$ and any admissible functions f_1, \dots, f_d such that their derivatives f'_1, \dots, f'_d are linearly independent modulo constants,⁹ the push-forward of μ under the mapping

$$u(\cdot) \mapsto \left(\int_{\mathbb{T}^2} f_k(\operatorname{rot} u(x)) \, dx, \, 1 \leq k \leq d \right) \in \mathbb{R}^d,$$

is absolutely continuous with respect to the Lebesgue measure.

Due to this result, the Hausdorff dimension of $\operatorname{supp} \mu$ is infinite. Indeed, if this is not the case, then choosing d bigger than the Hausdorff dimension, we see that the push-forward measure in item (2) of the theorem is supported by a set of dimension $< d$, which contradicts the assertion.

6 3d Navier–Stokes system in thin domains

In this section we consider a thin layer around the torus \mathbb{T}^2 and the 3d Navier–Stokes system, perturbed by a random kick-force with a sufficiently small vertical component. We show that when the width of the layer goes to zero, statistical characteristics of the 3d flow converge to those of the 2d flow (0.1), (1.2), (1.6), where the kicks η^k are the horizontal components of the 3d kicks. Moreover, this convergence is uniform in time. Since Earth’s atmosphere is a thin spherical layer, this result gives a good support to the belief that suitably chosen 2d stochastic meteorological models can be successfully used to model the climate. Usually these 2d models are related to Eq. (0.1); see the works [Var13, Kle17] and the references therein.

Let $Q_\varepsilon = \mathbb{T}^2 \times (0, \varepsilon) = \{x = (x_1, x_2, x_3)\}$. Consider the 3d Navier–Stokes system in Q_ε under the free boundary conditions:

$$\partial_t u + \langle u, \nabla \rangle u - \nu \Delta u + \nabla p = \eta(t, x), \quad \operatorname{div} u = 0, \quad x \in Q_\varepsilon, \quad (6.1)$$

$$u_3 = \partial_3 u_{1,2} = 0 \quad \text{for } x_3 = 0 \text{ and } x_3 = \varepsilon, \quad (6.2)$$

$$u(0, x) = u_0(x), \quad (6.3)$$

where $u = (u_1, u_2, u_3)$. Denote by H_ε (by V_ε) the L_2 -space (the H^1 -space) of divergence-free vector fields $(u_1, u_2, u_3)(x)$ on Q_ε such that u_1 and u_2 have zero mean. Denote by $|\cdot|_\varepsilon$ the L_2 -norm on Q_ε , i.e. the norm in H_ε (note that $|\mathbf{1}|_\varepsilon = \varepsilon$), by $(\cdot, \cdot)_\varepsilon$ the corresponding L_2 scalar product, and denote by $\|\cdot\|_\varepsilon$ the homogeneous norm in V_ε , $\|u\|_\varepsilon = |\nabla u|_\varepsilon$. The space H as in Section 1.1 is naturally embedded in H_ε :

$$i : H \ni (u_1, u_2) \mapsto (u_1(x_1, x_2), u_2(x_1, x_2), 0) \in H_\varepsilon,$$

and the norm of this embedding equals $\sqrt{\varepsilon}$. Introduce in H_ε two orthogonal projections:

$$M_\varepsilon u = \left(\varepsilon^{-1} \int_0^\varepsilon u_1(x', y) \, dy, \varepsilon^{-1} \int_0^\varepsilon u_2(x', y) \, dy, 0 \right), \quad N_\varepsilon = \operatorname{id} - M_\varepsilon,$$

⁹That is, if $c_1 f'_1 + \dots + c_d f'_d \equiv \text{const}$, then all c_k vanish.

where $x' = (x_1, x_2)$. Then $M_\varepsilon H_\varepsilon = iH$. The 3d Stokes operator $L_\varepsilon = -\Delta|_{H_\varepsilon}$ preserves the spaces $M_\varepsilon H_\varepsilon$ and $N_\varepsilon H_\varepsilon$, and its eigenfunctions are of two kinds:

$$e_s(x_1, x_2) \in M_\varepsilon H_\varepsilon, \quad s \in \mathbb{Z}_\varepsilon^2, \quad \text{and} \quad e_j^\varepsilon(x_1, x_2, x_3) \in N_\varepsilon H_\varepsilon, \quad j \in \mathbb{N},$$

where $|e_s|_\varepsilon = |e_j^\varepsilon|_\varepsilon = \sqrt{\varepsilon}$ for all s and j , and

$$L_\varepsilon e_s = |s|^2 e_s \quad \forall s, \quad L_\varepsilon e_j^\varepsilon = \Lambda_j^\varepsilon e_j^\varepsilon \quad \forall j,$$

so that

$$\|e_s\|_\varepsilon = (L e_s, e_s)_\varepsilon^{1/2} = |s| \sqrt{\varepsilon}, \quad \|e_j^\varepsilon\|_\varepsilon = \sqrt{\Lambda_j^\varepsilon} \varepsilon.$$

The vectors $e_s \in iH$ will be identified with the eigenvectors of the 2d Stokes operator (which is the opposite of the 2d Laplacian, restricted to the space H). Each vector e_j^ε has components $(e_j^\varepsilon)_l = C_j^l(\sin / \cos)(s_j^l \cdot x') \cos(\frac{\pi}{\varepsilon} n_j^l x_3)$, $l = 1, 2, 3$, where $s_j^l \in \mathbb{Z}_*^2$, $n_j^l \in \mathbb{N} \cup \{0\}$ and at least one of the numbers n_j^1, \dots, n_j^3 is non-zero. Therefore the eigenvalue Λ_j^ε has the form

$$\Lambda_j^\varepsilon = A_j + B_j \pi^2 \varepsilon^{-2}, \quad A_j \in \mathbb{N} \cup \{0\}, \quad B_j \in \mathbb{N}. \quad (6.4)$$

Assume that the force $\eta(t, x)$ in (6.1) is a kick-process of the form (1.6) with the kicks

$$\eta_\varepsilon^k(x) = \sum_{s \in \mathbb{Z}_*^2} b_s \xi_s^k e_s(x) + \sum_{j=1}^{\infty} d_j^\varepsilon \zeta_j^k e_j^\varepsilon(x).$$

Here the constants $\{b_s\}$ and $\{d_j^\varepsilon\}$ are such that

$$\mathfrak{B}_1 := \sum b_s^2 |s|^2 < \infty, \quad b_s \neq 0 \quad \forall s, \quad \mathfrak{D}_1^\varepsilon := \sum (d_j^\varepsilon)^2 \Lambda_j^\varepsilon < \infty, \quad (6.5)$$

and $\{\xi_s^k\}, \{\zeta_j^k\}$ are i.i.d. random variable with law $p(r)dr$, where $p \in C_0^1(-1, 1)$ satisfies the conditions

$$p(0) \neq 0, \quad \int_{\mathbb{R}} r p(r) dr = 0.$$

We shall compare solutions of Eqs. (6.1)–(6.3) with those of the kick-forced 2d Navier–Stokes system (0.1), (1.2), (1.6), where the kicks η^k are

$$\eta^k = M_\varepsilon \eta_\varepsilon^k = \sum_s b_s \zeta_s^k e_s(x).$$

Under the assumptions (6.5), this 2d equation is exponentially mixing in the space H in the sense that there exists a measure $\mu \in \mathcal{P}(H)$ such that for every solution $u(t)$ of (0.1), (1.2), (1.6) its law, evaluated at integer points $t \in \mathbb{N}$, converges to μ exponentially fast in the dual-Lipschitz norm in $\mathcal{P}(H)$; see [KS12]. The mixing property for Eqs. (6.1)–(6.3), claimed in the theorem below, is understood in a similar way.

Theorem 6.1. *Assume that*

$$\mathcal{D}_1^\varepsilon \leq \varepsilon^{-1}\gamma^2(\varepsilon), \quad \text{where } \gamma(\varepsilon) \rightarrow 0 \text{ as } \varepsilon \rightarrow 0. \quad (6.6)$$

Then there exist $c_0, \varepsilon_0 > 0$ such that the following properties hold.

Exponential mixing. *If $0 < \varepsilon \leq \varepsilon_0$, then the set*

$$\mathcal{O}_\varepsilon = \{u : \|M_\varepsilon u\|_\varepsilon \leq c_0\sqrt{\varepsilon}, \|N_\varepsilon u\|_\varepsilon \leq c_0\gamma(\varepsilon)\} \subset V_\varepsilon$$

is invariant for Eqs. (6.1)–(6.3) and the dynamics on \mathcal{O}_ε is exponentially mixing with invariant measure $\mu_\varepsilon \in \mathcal{P}(\mathcal{O}_\varepsilon)$.

Convergence. *As $\varepsilon \rightarrow 0$, the measure $(M_\varepsilon)_*\mu_\varepsilon$ converges to μ weakly in H .*

Stability. *Let $v_0 \in H$ be such that $\|v_0\|_1 < c_0$, let $u(t)$ be a solution of (0.1), (1.2), (1.6), equal v_0 at $t = 0$, and let $u_\varepsilon(t)$ be a solution of Eqs. (6.1)–(6.3) with $u_0 = iv_0 \in \mathcal{O}_\varepsilon$. Then*

$$\|\mathcal{D}(M_\varepsilon u_\varepsilon(t)) - \mathcal{D}(u(t))\|_{L(H)}^* \rightarrow 0 \quad \text{as } \varepsilon \rightarrow 0,$$

uniformly in $t \geq 0$.

For a random vector field u on Q_ε its averaged normalised kinetic energy is $\mathcal{E}_\varepsilon(u) = \frac{1}{2}\mathbb{E}|u|_\varepsilon^2 / \text{Vol}(Q_\varepsilon) = \frac{1}{2\varepsilon}\mathbb{E}|u|_\varepsilon^2$. By (6.6) and (6.4), for a kick η_ε^k , we have

$$\mathcal{E}_\varepsilon(M_\varepsilon \eta_\varepsilon^k) = \frac{1}{2}\kappa^2 \sum_s b_s^2 \sim 1, \quad \mathcal{E}_\varepsilon(N_\varepsilon \eta_\varepsilon^k) = \frac{1}{2}\kappa^2 \sum_j (d_j^\varepsilon)^2 \lesssim \varepsilon\gamma^2(\varepsilon)$$

where $\kappa^2 = \int r^2 p(r) dr$. For the averaged normalised dissipation of energy $\mathcal{D}_\varepsilon(u) = \frac{1}{2\varepsilon}\mathbb{E}\|u\|_\varepsilon^2$, we have

$$\mathcal{D}_\varepsilon(M_\varepsilon \eta_\varepsilon^k) = \frac{1}{2}\kappa^2 \mathfrak{B}_1 \sim 1, \quad \mathcal{D}_\varepsilon(N_\varepsilon \eta_\varepsilon^k) = \frac{1}{2}\kappa^2 \mathcal{D}_1^\varepsilon \leq \frac{1}{2}\kappa^2 \varepsilon^{-1}\gamma^2(\varepsilon).$$

Therefore the vertical component of the random force in (6.1) should be small in terms of the energy, but not in terms of the dissipation of energy.

7 Open problems

Open problem 1 (Mixing of pipe flow). Let us consider the Navier–Stokes system (0.1) in the strip

$$D = \{(x_1, x_2) \in \mathbb{R}^2 : x_1 \in \mathbb{R}, |x_2| < 1\}.$$

The initial-boundary value problem for (0.1) is well posed in L^∞ spaces on D , with no decay conditions at infinity (see [AM05, AZ14]). Moreover, the problem is dissipative and possesses a global attractor (in the deterministic setting). It follows that, at least in the case of a bounded stochastic forcing, the random dynamics has a stationary distribution. A challenging problem is to prove its *uniqueness* and *mixing*.

Open problem 2 (Mixing in the whole space). The 2d Navier–Stokes system considered on the whole space \mathbb{R}^2 is well posed in L^∞ spaces (see [GMS01, LR02]). However, it is not known if the dynamics is dissipative, and it seems to be a hopeless task to prove any kind of regular behaviour of solutions under stochastic perturbations. On the other hand, the Navier–Stokes system with the Ekman damping

$$\partial_t u + \gamma u + \langle u, \nabla \rangle u - \nu \Delta u + \nabla p = f(t, x), \quad \operatorname{div} u = 0, \quad (7.1)$$

where $\gamma > 0$ is a number, is dissipative (see [Zel13]). It is a natural question to investigate the *existence of stationary measure* and its *mixing properties* under various types of random perturbation.

Open problem 3 (Vanishing noise limit). The vanishing noise limit of stationary measures described in Theorem 4.2 concerns a rather particular situation: the unperturbed dynamics should be globally asymptotically stable. This is a very restrictive hypothesis, and removing it is an important question. In the finite-dimensional case this problem is rather well understood, and one can establish the so-called *Freidlin–Wentzell asymptotics* for stationary measures; see Section 6.4 in [FW12]. As for stochastic PDE’s, similar results can be proved, provided that the global attractor for the unperturbed dynamics has a regular structure. The latter means that the attractor consists of finitely many steady states and the heteroclinic orbits joining them. Such a result is proved in [Mar17b] for the case of a damped nonlinear wave equation. The attractor of the Navier–Stokes system is not likely to possess that property, and the validity of the *Freidlin–Wentzell type asymptotics* for stationary distributions remains an open problem.

Open problem 4 (Inviscid limit). (a) Does the limiting measure μ in (5.6) depend on the sequence $\{\mu_{\nu_j}\}$? (We believe that it does).

(b) Consider the energy spectrum of a measure μ as above:

$$E_k(\mu) = Z^{-1} \sum_{\{s \in \mathbb{Z}_+^2 : M^{-1}k \leq |s| \leq Mk\}} \mathbb{E}^\mu |u_s|^2,$$

where Z is the number of terms in the sum and $M > 1$ is a suitable constant. Do there exist positive constants a, b, C such that

$$C^{-1}k^{-a} \leq E_k(\mu) \leq Ck^{-b} \quad (7.2)$$

for all $k \geq 1$? Do the exponents a and b depend on the inviscid limit μ ?

(c) Do the stationary measures μ^ν satisfy (7.2) for $C_1\nu^{-\alpha} \leq k \leq C_2\nu^{-\beta}$, for suitable $0 \leq \alpha < \beta$ and $C_1, C_2 > 0$? See [Bor13] for an affirmative answer to this question with $a = b = 2$ and $\alpha = 0, \beta = 1$, when u is a solution of the 1d stochastic Burgers equation.

Open problem 5 (thin 3d domains). Improve the result of Theorem 6.1 by replacing condition (6.6) with a weaker restraint, thus allowing in (6.1) for random forces with bigger vertical components (this seems to be possible to achieve

by making better use of the stochastic nature of the force η). Obtain similar results for 3d stochastic models of Earth's atmosphere and their suitable 2d approximations.

Frequently used notations

C, C_i	positive numbers that may depend on the parameters mentioned in brackets
$J \subset \mathbb{R}$	closed interval
$C(J, E)$	space of continuous functions on J with range in a Banach space E
$L^p(J, E)$	space of Borel-measurable functions $f : J \rightarrow E$ such that $\int_J \ f(t)\ _E^p dt < \infty$, with obvious modification for $p = \infty$
X	complete separable metric space with a distance d_X
$\mathcal{B}(X)$	Borel σ -algebra of X
$C_b(X)$	space of bounded continuous functions $g : X \rightarrow \mathbb{R}$ with the norm $\ g\ _\infty = \sup_{u \in X} g(u) $
$d_X(u, B)$	minimal distance from a point $u \in X$ to a subset $B \subset X$
$\mathcal{P}(X)$	space of Borel probability measures on X endowed with the topology of weak convergence; the latter can be metrised by the Lipschitz-dual metric $\ \cdot\ _{\text{Lip}(X)}^*$, see (2.8)
$\mathcal{D}(\xi)$	the law of a random variable ξ
$F_*\mu$	the image of a measure μ under a measurable mapping F
$L^p(D), H^s(D)$	Lebesgue and Sobolev spaces on a domain $D \subset \mathbb{R}^d$ with standard norms $\ \cdot\ _p$ and $\ \cdot\ _s$, respectively; sometimes we write $L^p(D, \mathbb{R}^m)$ and $H^s(D, \mathbb{R}^m)$ to denote the corresponding spaces of \mathbb{R}^m -valued functions, and (\cdot, \cdot) stands for L_2 -scalar product
H^s	space of divergence-free vector functions with zero mean that belong to $H^s(\mathbb{T}^2, \mathbb{R}^2)$, $s \in \mathbb{N} \cup \{0\}$.
H, V	abbreviations for the spaces H^0 and H^1
L_σ^2	the set of divergence-free functions in $L^2(\mathbb{T}^2, \mathbb{R}^2)$
$\mathbb{R}_+, \mathbb{Z}_+$	sets of real numbers and non-negative integers, respectively
$\mathbb{T}^2 = \mathbb{R}^2/2\pi\mathbb{Z}^2$	two-dimensional torus with sides 2π
$(\Omega, \mathcal{F}, \mathbb{P})$	complete probability space

References

- [AM05] A. L. Afendikov and A. Mielke, *Dynamical properties of spatially non-decaying 2D Navier-Stokes flows with Kolmogorov forcing in an infinite strip*, J. Math. Fluid Mech. **7** (2005), no. suppl. 1, S51–S67.
- [AZ14] P. Anthony and S. Zelik, *Infinite-energy solutions for the Navier-Stokes equations in a strip revisited*, Commun. Pure Appl. Anal. **13** (2014), no. 4, 1361–1393.

- [Bat82] G. K. Batchelor, *The Theory of Homogeneous Turbulence*, Cambridge University Press, Cambridge, 1982.
- [BC17] Z. Brzeźniak and S. Cerrai, *Large deviation principle for the invariant measures of the 2D stochastic Navier–Stokes equations on a torus*, J. Funct. Anal. **273** (2017), no. 6, 1891 – 1930.
- [BKL02] J. Bricmont, A. Kupiainen, and R. Lefevere, *Exponential mixing of the 2D stochastic Navier–Stokes dynamics*, Comm. Math. Phys. **230** (2002), no. 1, 87–132.
- [Bor13] A. Boritchev, *Sharp estimates for turbulence in white-forced generalised Burgers equation*, Geom. Funct. Anal. **23** (2013), no. 6, 1730–1771.
- [BV12] F. Bouchet and A. Venaille, *Statistical mechanics of two-dimensional and geophysical flows*, Phys. Rep. **515** (2012), no. 5, 227–295.
- [CM10] I. Chueshov and A. Millet, *Stochastic 2D hydrodynamical type systems: Wellposedness and large deviations*, Appl. Math. Optim. **61** (2010), no. 3, 379–420.
- [Cra91] H. Crauel, *Markov measures for random dynamical systems*, Stochastics Stochastics Rep. **37** (1991), no. 3, 153–173.
- [DV75] M. D. Donsker and S. R. S. Varadhan, *Asymptotic evaluation of certain Markov process expectations for large time, I-II*, Comm. Pure Appl. Math. **28** (1975), 1–47, 279–301.
- [EMS01] W. E, J. C. Mattingly, and Ya. Sinai, *Gibbsian dynamics and ergodicity for the stochastically forced Navier–Stokes equation*, Comm. Math. Phys. **224** (2001), no. 1, 83–106.
- [Fla94] F. Flandoli, *Dissipativity and invariant measures for stochastic Navier–Stokes equations*, NoDEA Nonlinear Differential Equations Appl. **1** (1994), no. 4, 403–423.
- [FM95] F. Flandoli and B. Maslowski, *Ergodicity of the 2D Navier–Stokes equation under random perturbations*, Comm. Math. Phys. **172** (1995), no. 1, 119–141.
- [FP67] C. Foias and G. Prodi, *Sur le comportement global des solutions non-stationnaires des équations de Navier–Stokes en dimension 2*, Rend. Sem. Mat. Univ. Padova **39** (1967), 1–34.
- [Fri95] U. Frisch, *Turbulence. The Legacy of A. N. Kolmogorov*, Cambridge University Press, Cambridge, 1995.
- [FW12] M. I. Freidlin and A. D. Wentzell, *Random Perturbations of Dynamical Systems*, Springer, Heidelberg, 2012.

- [Gal02] G. Gallavotti, *Foundations of Fluid Dynamics*, Springer-Verlag, Berlin, 2002.
- [GMS01] Y. Giga, S. Matsui, and O. Sawada, *Global existence of two-dimensional Navier–Stokes flow with nondecaying initial velocity*, J. Math. Fluid Mech. **3** (2001), no. 3, 302–315.
- [GSV15] N. Glatt-Holtz, V. Sverak, and V. Vicol, *On inviscid limits for the stochastic Navier–Stokes equations and related models*, Arch. Rat. Mech. Anal. **217** (2015), no. 1, 619–649.
- [HM06] M. Hairer and J. C. Mattingly, *Ergodicity of the 2D Navier–Stokes equations with degenerate stochastic forcing*, Ann. of Math. (2) **164** (2006), no. 3, 993–1032.
- [HM11] ———, *A theory of hypoellipticity and unique ergodicity for semilinear stochastic PDEs*, Electron. J. Probab. **16** (2011), no. 23, 658–738.
- [JNPS15] V. Jakšić, V. Nersesyan, C.-A. Pillet, and A. Shirikyan, *Large deviations from a stationary measure for a class of dissipative PDEs with random kicks*, Comm. Pure Appl. Math. **68** (2015), no. 12, 2108–2143.
- [JNPS17] ———, *Large deviations and mixing for dissipative PDE’s with unbounded random kicks*, Nonlinearity (2017), to appear.
- [Kle17] Yu. Yu. Klevtsova, *On the rate of convergence of distributions of solutions to the stationary measure as $t \rightarrow +\infty$ for the stochastic system of the Lorenz model describing a baroclinic atmosphere*, Mat. Sb. **208** (2017), no. 7, 19–67.
- [KNS18] S. Kuksin, V. Nersesyan, and A. Shirikyan, *Exponential mixing for a class of dissipative PDEs with bounded degenerate noise*, in preparation (2018).
- [KP05] S. Kuksin and O. Penrose, *A family of balance relations for the two-dimensional Navier–Stokes equations with random forcing*, J. Statist. Phys. **118** (2005), no. 3-4, 437–449.
- [KS91] I. Karatzas and S. E. Shreve, *Brownian Motion and Stochastic Calculus*, Springer-Verlag, New York, 1991.
- [KS00] S. Kuksin and A. Shirikyan, *Stochastic dissipative PDEs and Gibbs measures*, Comm. Math. Phys. **213** (2000), no. 2, 291–330.
- [KS01] ———, *A coupling approach to randomly forced nonlinear PDE’s. I*, Comm. Math. Phys. **221** (2001), no. 2, 351–366.
- [KS02] ———, *Coupling approach to white-forced nonlinear PDEs*, J. Math. Pures Appl. (9) **81** (2002), no. 6, 567–602.

- [KS03] ———, *Some limiting properties of randomly forced two-dimensional Navier–Stokes equations*, Proc. Roy. Soc. Edinburgh Sect. A **133** (2003), no. 4, 875–891.
- [KS04] ———, *On random attractors for systems of mixing type*, Funktsional. Anal. i Prilozhen. **38** (2004), no. 1, 34–46, 95.
- [KS12] ———, *Mathematics of Two-Dimensional Turbulence*, Cambridge University Press, Cambridge, 2012.
- [Kuk02a] S. Kuksin, *Ergodic theorems for 2D statistical hydrodynamics*, Rev. Math. Phys. **14** (2002), no. 6, 585–600.
- [Kuk02b] ———, *On exponential convergence to a stationary measure for non-linear PDEs perturbed by random kick-forces, and the turbulence limit*, Partial differential equations, Amer. Math. Soc. Transl. Ser. 2, vol. 206, Amer. Math. Soc., Providence, RI, 2002, pp. 161–176.
- [Kuk04] ———, *The Eulerian limit for 2D statistical hydrodynamics*, J. Statist. Phys. **115** (2004), no. 1-2, 469–492.
- [Kuk08] ———, *On distribution of energy and vorticity for solutions of 2D Navier–Stokes equation with small viscosity*, Comm. Math. Phys. **284** (2008), no. 2, 407–424.
- [Lam96] J. Lamperti, *Probability*, John Wiley & Sons, New York, 1996.
- [Le 87] Y. Le Jan, *Équilibre statistique pour les produits de difféomorphismes aléatoires indépendants*, Ann. Inst. H. Poincaré Probab. Statist. **23** (1987), no. 1, 111–120.
- [Led86] F. Ledrappier, *Positivity of the exponent for stationary sequences of matrices*, Lyapunov Exponents (Bremen, 1984), Springer, Berlin, 1986, pp. 56–73.
- [LR02] P.-G. Lemarié-Rieusset, *Recent Developments in the Navier–Stokes Problem*, Chapman & Hall/CRC, Boca Raton, FL, 2002.
- [Mar17a] D. Martirosyan, *Large deviations for invariant measures of white-forced 2D Navier–Stokes equation*, J. Evol. Equ. (2017), accepted for publication.
- [Mar17b] ———, *Large deviations for stationary measures of stochastic non-linear wave equations with smooth white noise*, Comm. Pure Appl. Math. **70** (2017), no. 9, 1754–1797.
- [Mat02] J. C. Mattingly, *Exponential convergence for the stochastically forced Navier–Stokes equations and other partially dissipative dynamics*, Comm. Math. Phys. **230** (2002), no. 3, 421–462.

- [MN17] D. Martirosyan and V. Nersesyan, *Local large deviations principle for occupation measures of the damped nonlinear wave equation perturbed by a white noise*, Ann. Inst. H. Poincaré Probab. Statist. (2017), to appear.
- [Ner17] V. Nersesyan, *Large deviations for the Navier–Stokes equations driven by a white-in-time noise*, in preparation (2017).
- [Oda08] C. Odasso, *Exponential mixing for stochastic PDEs: the non-additive case*, Probab. Theory Related Fields **140** (2008), no. 1-2, 41–82.
- [Shi05] A. Shirikyan, *Ergodicity for a class of Markov processes and applications to randomly forced PDE’s. I*, Russ. J. Math. Phys. **12** (2005), no. 1, 81–96.
- [Shi06] ———, *Law of large numbers and central limit theorem for randomly forced PDE’s*, Probab. Theory Related Fields **134** (2006), no. 2, 215–247.
- [Shi15] ———, *Control and mixing for 2D Navier-Stokes equations with space-time localised noise*, Ann. Sci. Éc. Norm. Supér. (4) **48** (2015), no. 2, 253–280.
- [Shi18] ———, *Controllability implies mixing II. Convergence in the dual-Lipschitz metric*, in preparation (2018).
- [Tem79] R. Temam, *Navier–Stokes Equations*, North-Holland, Amsterdam, 1979.
- [Var13] G. A. Varner, *Stochastically Perturbed Navier–Stokes System on the Rotating Sphere*, ProQuest LLC, Ann Arbor, MI, 2013.
- [VF88] M. I. Vishik and A. V. Fursikov, *Mathematical Problems in Statistical Hydromechanics*, Kluwer, Dordrecht, 1988.
- [Zel13] S. Zelik, *Infinite energy solutions for damped Navier-Stokes equations in \mathbb{R}^2* , J. Math. Fluid Mech. **15** (2013), no. 4, 717–745.