# Remarks on the Balance Relations for the Two-Dimensional Navier-Stokes Equation with Random Forcing 

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#### Abstract

We use the balance relations for the stationary in time solutions of the randomly forced 2D Navier-Stokes equations, found in [10], to study these solutions further. We show that the vorticity $\xi(t, x)$ of a stationary solution has a finite exponential moment, and that for any $a \in \mathbb{R}, t \geq 0$ the expectation of the integral of $\left|\nabla_{x} \xi\right|$ over the level-set $\{x \mid \xi(t, x)=a\}$, up to a constant factor equals the expectation of the integral of $\left|\nabla_{x} \xi\right|^{-1}$ over the same set.


KEY WORDS: Two-dimensional Navier-Stokes equation; stationary measure; vorticity; balance relations; exponential moment.

## 1. INTRODUCTION

In this paper we continue to study the 2D Navier-Stokes equation on the two-dimensional torus $\mathbb{T}^{2}=\mathbb{R}^{2} /\left(2 \pi \mathbb{Z}^{2}\right)$, perturbed by a random force:

$$
\begin{array}{r}
\dot{u}-v \Delta u+(u \cdot \nabla) u+\nabla p(t, x)=\sqrt{v} \tilde{\eta}(t, x) \\
u=u(t, x) \in \mathbb{R}^{2}, \quad x \in \mathbb{T}^{2}, \quad \operatorname{div} u=0, \quad \int u d x \equiv \int \tilde{\eta} d x \equiv 0 \tag{1.1}
\end{array}
$$

The scaling factor $\sqrt{v}$ in the r.h.s. of the equation is not important for us (since the force $\tilde{\eta}$ may depend on $\nu$ ), but it makes the formulas obtained nicer.

The force $\tilde{\eta}$ is a stationary random field in $t$ and $x$, smooth periodic in $x$ and white in $t$. If $\tilde{\eta}$ satisfies some mild nondegeneracy assumptions, discussed in the next section, then the probability distribution of any solution $u(t, x)$ of the stochastic differential equation (1.1) converges as $t \rightarrow \infty$ to a unique stationary measure (which is a measure on a function space, forms by divergence-free vector

[^0]fields $u(x))$. Let $u_{v}(t, x), t \geq 0$, be a stationary in time solution of (1.1), for any $t$ distributed as the stationary measure, and let $\xi_{v}(t, x)$ be its vorticity:
$$
\xi_{v}(t, x)=\operatorname{curl} u_{v}=\frac{\partial u_{v}^{2}}{\partial x_{1}}-\frac{\partial u_{v}^{1}}{\partial x_{2}}
$$

The vorticity $\xi_{v}$ satisfies the diffusion-convection equation

$$
\begin{equation*}
\dot{\xi}-v \Delta \xi+(u \cdot \nabla) \xi=\sqrt{v} \operatorname{curl} \tilde{\eta} \tag{1.2}
\end{equation*}
$$

Since $\tilde{\eta}$ is stationary in $x$, then both $u_{v}$ and $\xi_{v}$ are stationary in $x$ as well.
It is established in [10] that the process $\xi_{v}$ satisfies infinitely many balance relations:

$$
\begin{equation*}
\mathbf{E}\left(g\left(\xi_{v}(t, x)\right)\left|\nabla \xi_{v}(t, x)\right|^{2}\right)=B_{1} \mathbf{E}\left(g\left(\xi_{v}(t, x)\right) \quad \forall t, x\right. \tag{1.3}
\end{equation*}
$$

Here $g$ is any continuous function which has at most a polynomial growth, i.e. $|g(v)| \leq C\left(1+|v|^{k}\right)$ for all $v$, where $C$ and $k$ are some fixed constants. The constant $B_{1}$ is explicitly defined in terms of the force $\tilde{\eta}$. The relations (1.3) are related to the Helmholtz invariants for inviscid 2d flow.

The goal of this work is to derive from the balance relations two corollaries. The first corollary ${ }^{2}$ deals with integrals over level-sets $\Gamma(\tau)$ of the vorticity $\xi_{\nu}$,

$$
\Gamma(\tau)=\left\{x \mid \xi_{v}(t, x)=\tau\right\}
$$

where $t \geq 0$ is fixed. Each $\Gamma(\tau)$ is a random subset of the torus $\mathbb{T}^{2}$ (the notation of a random parameter is suppressed everywhere in the Introduction). In Theorem (3.2), Section (3), we derive the following co-area form of the balance relations:

$$
\begin{equation*}
\mathbf{E} \int_{\Gamma(\tau)}\left|\nabla \xi_{v}\right| d \gamma=B_{1} \mathbf{E} \int_{\Gamma(\tau)}\left|\nabla \xi_{v}\right|^{-1} d \gamma \tag{1.4}
\end{equation*}
$$

for almost all $\tau$. Here $d \gamma$ is the length element on the random curve $\Gamma(\tau)$, which is well defined a.s.

If $\tau$ is a regular value of $\xi_{v}$ as a function of $x$, then $\Gamma(\tau)$ is a finite union of smooth curves, and in any point $x \in \Gamma(\tau)$ we have $\nabla \xi_{v}= \pm \frac{\partial}{\partial n} \xi_{\nu}$, where $n$ is a unit normal to the curve. Hence, in (1.4) $\left|\nabla \xi_{\nu}\right|$ can be replaced by $\left|\frac{\partial}{\partial n} \xi_{\nu}\right|$. So the integral $\int\left|\nabla \xi_{v}\right| d \gamma$ is a sum of the moduli of the flows of $\nabla \xi_{v}$ through the curves, forming the set $\Gamma$. The integral $\int\left|\nabla \xi_{\nu}\right|^{-1} d \gamma$ can be interpreted similar.

If $u(t, x)$ is any solution of $(1.1)$ and $\xi(t, x)$ is its vorticity, then by the ergodic theorem the average in ensemble can be replaced by the average in time. Therefore, $\xi(t, x)$ satisfies (1.4), where the expectation $\mathbf{E}$ is replaced by $\lim T^{-1} \int_{1}^{T+1} \ldots d t$ (cf. the end of Section (2), where the balance relations (1.3) are re-interpreted similar).

[^1]The level sets $\Gamma(\tau)$ of the vorticity of a solution for the deterministic NavierStokes equation in the 2 d and 3d cases, and of solutions for equation (1.2) without assuming that $\xi=\operatorname{curl} u$ (but imposing certain a-priori bounds on $u$ and $\xi$ ), were studied by P. Constantin and others in $[2,4,5]$. There the areas of the sets $\Gamma(\tau)$ are estimated (with and without averaging in $t$ and $\tau$ ), as well as certain integrals over these sets. The results of $[2,5]$ are physically motivated. At this moment we cannot suggest any physical interpretation of the relations (1.4). We simply believe that they are important as exact relations, satisfied by solutions of a basic equation of mathematical physics.

In Section (4) the balance relations are used to prove that the random field $\xi_{v}$ has finite exponential moments:

$$
\mathbf{E} e^{\sigma\left|\xi_{v}(t, x)\right|} \leq C<\infty
$$

for some $\sigma>0$, uniformly in $t, x$ and in $v>0$. Moreover, $\mathbf{E} e^{\sigma_{1}\left|u_{v}(t, x)\right|} \leq C_{1}$ and $\mathbf{E} e^{\sigma_{2}\left|\nabla_{x} u_{v}(t, x)\right|^{1 / 2}} \leq C_{2}$ (for all $t, x$, uniformly in $v>0$ ). In particular,

$$
\mathbf{P}\left\{\left|\xi_{v}(t, x)\right| \geq K\right\} \leq C e^{-\sigma K} \quad \forall K
$$

etc. That is, high values of $u_{\nu}(t, x)$, or of $\xi_{v}(t, x)$, or of $\nabla u_{v}(t, x)$ are very unlikely.

## 2. PRELIMINARIES

We denote by $H$ the space of square-integrable vector fields $u(x)$ such that $\operatorname{div} u=0$ and $\int u d x=0$ given the $L_{2}$-norm $\|\cdot\|$ and the $L_{2}$-scalar product $(\cdot, \cdot)$. By $H^{n}, n \in \mathbb{N}$, we denote the Sobolev space $H^{n}=H \cap H^{n}\left(\mathbb{T}^{2} ; \mathbb{R}^{2}\right)$ given the norm $\|u\|_{n}=\left((-\Delta)^{n} u, u\right)^{1 / 2}$.

Let us denote by $\Pi$ the Leray projector $\Pi: L_{2}\left(\mathbb{T}^{2} ; \mathbb{R}^{2}\right) \rightarrow H$, which removes the gradient and the constant part of a vector field it operates upon (see e.g., [3]). Applying $\Pi$ to (1.1) we write it in the usual form

$$
\begin{equation*}
\dot{u}(t, x)+v L u+B(u)=\sqrt{v} \eta, \tag{2.1}
\end{equation*}
$$

where we have denoted $L u=-\Pi \Delta u, B(u)=\Pi(u \cdot \nabla) u$ and $\eta=\Pi \tilde{\eta}$.
The force $\eta(t, x)$ is assumed to be a Gaussian random field, white in time and smooth in $x$ :

$$
\begin{equation*}
\eta=\frac{d}{d t} \zeta(t, x), \quad \zeta=\sum_{s \in \mathbb{Z}_{0}^{2}} b_{s} \beta_{s}(t) e_{s}(x) \tag{2.2}
\end{equation*}
$$

Here $\mathbb{Z}_{0}^{2}=\mathbb{Z}^{2} \backslash\{0\},\left\{\beta_{s}(t)=\beta_{s}^{\omega}(t)\right\}$ is a set of independent standard Wiener processes, defined on a complete probability space $(\Omega, \mathcal{F}, \mathbf{P})$, satisfying $\beta_{s}(0)=0$.

The real coefficients $b_{s}$ are such that $M_{1}<\infty$, where

$$
M_{n}=\sum_{s \in \mathbb{Z}_{0}^{2}}|s|^{2 n} b_{s}^{2}, \quad n \in \mathbb{Z}
$$

and $\left(e_{s}, s \in \mathbb{Z}_{0}^{2}\right)$ is the standard trigonometric basis of $H$ :

$$
\begin{array}{ll}
e_{s}(x)=\frac{\sin (s \cdot x)}{\sqrt{2} \pi|s|}\left[\begin{array}{c}
-s_{2} \\
s_{1}
\end{array}\right] & \text { if } s_{1}+s_{2} \delta_{s_{1}, 0}>0 \\
e_{s}(x)=\frac{\cos (s \cdot x)}{\sqrt{2} \pi|s|}\left[\begin{array}{c}
-s_{2} \\
s_{1}
\end{array}\right] & \text { if } s_{1}+s_{2} \delta_{s_{1}, 0}<0
\end{array}
$$

The vector fields $e_{s}$ are eigenvectors of $L: L e_{s}=|s|^{2} e_{s}$ for each $s$.
We shall often interpret random fields $u(t, x)$ and $\eta(t, x)$ as random processes in $H$ (or in another space $H^{n}$ ) and write them as $u(t)$ and $\eta(t)$.

It is well known (see $[6,8,18]$ ) that for any random initial data $u_{0}$, independent of $\eta(\cdot)$, and such that $\mathbf{E}\left\|u_{0}\right\|^{2}<\infty$, the equation (2.1) has a unique solution $u(t)$, belonging to $C([0, \infty), H)$ almost sure. This solution is a Markov process in $H$. We shall denote it as $u=u\left(t ; u_{0}\right)$. If

$$
\begin{equation*}
b_{s} \neq 0 \quad \forall s \tag{2.3}
\end{equation*}
$$

then the Markov process has a unique stationary measure $\mu_{\nu},{ }^{3}$ and all solutions $u(t)$ as above converge to $\mu_{\nu}$ in distribution:

$$
\mathcal{D} u(t) \rightharpoonup \mu_{v} \quad \text { as } \quad t \rightarrow \infty .
$$

Here and below $\mathcal{D}$ signifies the distribution of a random variable. If we choose for $u_{0}$ a random vector in $H$, distributed as $\mu_{v}$, then the corresponding solution, which will be denoted as $u_{v}(t), t \geq 0$, is stationary:

$$
\mathcal{D}\left(u_{v}(t)\right)=\mu_{v} \quad \forall t \geq 0
$$

Now let us assume that, in addition to (2.3), the coefficients $b_{s}$ are symmetric in $s$ :

$$
\begin{equation*}
b_{s}=b_{-s} \neq 0 \quad \forall s \tag{2.4}
\end{equation*}
$$

Then the random field $\zeta(t, x)$ is translationary invariant. That is, its distribution is invariant under the translation $T_{h}$ of the torus $\mathbb{T}^{2}, T_{h} x=x+h\left(h \in \mathbb{T}^{2}\right)$. Due to the uniqueness, this implies that the stationary measure $\mu_{\nu}$ and the stationary

[^2]process $u_{\nu}(t, x)$ also are translationary invariant:
$$
T_{h} \mu_{v}=\mu_{v}, \quad \mathcal{D} u_{\nu}(\cdot, \cdot+h)=\mathcal{D} u_{\nu}(\cdot, \cdot) \quad \forall h \in \mathbb{T}^{2}
$$
see [10, 14]. Under assumptions (2.4) the vorticity $\xi_{\nu}=\operatorname{curl} u_{\nu}$ satisfies the balance relations (1.3), where $B_{1}=1 / 2(2 \pi)^{-2} M_{1}$, see [10].

It is known (e.g., see $[18,8,13])$ that if $M_{k}<\infty(k \geq 1)$, and $u_{0} \in H$ is a non-random vector, then for any $T \geq 2$ the solution $u=u\left(t ; u_{0}\right)$ satisfies

$$
\begin{equation*}
u \in C\left([1, T], H^{k}\right) \quad \text { and } \quad u-\sqrt{v} \zeta \in C^{1}\left([1, T], H^{k-2}\right) \quad \text { a.s. . } \tag{2.5}
\end{equation*}
$$

The stationary solutions $u_{v}(t, x)$ also satisfy (2.5). Moreover, in this case the relations hold with $[1, T]$ replaced by any finite segment $\left[T_{1}, T_{2}\right], T_{1} \geq 0$.

For $u_{0} \in H$, let us denote $\xi\left(t, x ; u_{0}\right)=\operatorname{curl} u\left(t, x ; u_{0}\right)$. Noting that for any $x \in \mathbb{T}^{2}$ the map $H^{4} \rightarrow \mathbb{R}^{2}, u(\cdot) \mapsto\left(\operatorname{curl} u(x),|\nabla \operatorname{curl} u(x)|^{2}\right)$, is continuous, we see that

$$
\frac{1}{T} \int_{1}^{T+1} g\left(\xi\left(t, x ; u_{0}\right)\right)\left|\nabla \xi\left(t, x ; u_{0}\right)\right|^{2} d t \rightarrow \mathbf{E}\left(g\left(\xi_{v}(x)\left|\nabla \xi_{v}(x)\right|^{2}\right) \quad\right. \text { a.s. }
$$

by the Strong Law of Large Numbers (see [14, 16]), if

$$
\begin{equation*}
M_{5}<\infty, \quad \mathbf{E}\left|u_{0}\right|^{2}<\infty, \quad g(\xi) \text { is a polynomial. } \tag{2.6}
\end{equation*}
$$

Similar result holds for the functional $\xi \rightarrow g(\xi(t, x))$, so the balance relation (1.3) still holds if we replace $\mathbf{E}$ by $\lim T^{-1} \int_{1}^{T+1} \ldots d t$ and replace $\xi_{v}$ by any $\xi\left(t, x ; u_{0}\right)$, provided that (2.6) is satisfied.

## 3. THE CO-AREA FORM OF THE BALANCE RELATIONS

From now on we assume that

$$
M_{6}<\infty
$$

Then, by (2.5), there exists a null-set $\Omega_{0}$ such that for $\omega \notin \Omega_{0}$ we have

$$
\begin{equation*}
\xi_{v} \in C\left([0, T] ; C^{3}\left(\mathbb{T}^{2}\right)\right), \quad \xi_{v}-\sqrt{v} \operatorname{curl} \zeta \in C^{1}\left([0, T] \times \mathbb{T}^{2}\right) \tag{3.1}
\end{equation*}
$$

We re-define $\xi_{v}$ and $\zeta$ to vanish for $\omega \in \Omega_{0}$. Now (3.1) holds for all $\omega$.
Lemma 3.1. For any $t, x$ and $v$ we have

$$
\begin{equation*}
\mathbf{P}\left\{\nabla \xi_{\nu}(t, x)=0\right\}=0 \tag{3.2}
\end{equation*}
$$

The lemma is proved in Appendix.
Let us fix any $t \geq 0$, abbreviate $\xi_{v}(t, x)=\xi(x)$ and integrate the balance relation over $\mathbb{T}^{2}$ :

$$
\begin{equation*}
\mathbf{E} \int_{\mathbb{T}^{2}} g\left(\xi(x)|\nabla \xi(x)|^{2} d x=B_{1} \mathbf{E} \int_{\mathbb{T}^{2}} g(\xi(x)) d x\right. \tag{3.3}
\end{equation*}
$$

For a continuous function $g$ such that $|g| \leq 1$, we denote

$$
I_{1}=\mathbf{E} \int_{\mathbb{T}^{2}} g(\xi(x))|\nabla \xi(x)|^{2} d x, \quad I_{2}=\mathbf{E} \int_{\mathbb{T}^{2}} g(\xi(x)) d x
$$

Then

$$
\begin{equation*}
I_{1}=B_{1} I_{2} \tag{3.4}
\end{equation*}
$$

For $\varepsilon>0$ let us consider the set $K_{\varepsilon} \subset \mathbb{T}^{2} \times \Omega$,

$$
K_{\varepsilon}=\{(x, \omega)| | \nabla \xi \mid \geq \varepsilon\}
$$

Clearly,

$$
I_{1}=I_{1}^{\varepsilon}+O\left(\varepsilon^{2}\right), \quad I_{1}^{\varepsilon}=\mathbf{E} \int_{\mathbb{T}^{2}} g\left(\xi(x)|\nabla \xi(x)|^{2} I_{K_{\varepsilon}}(x, \omega) d x\right.
$$

For the quantity $I_{2}^{\varepsilon}$, obtained from $I_{2}$ by multiplying the integrand by the factor $I_{K_{\varepsilon}}$, we have

$$
\left|I_{2}-I_{2}^{\varepsilon}\right| \leq \int_{\Omega} \int_{T^{2}} I_{\{|\nabla \xi|<\varepsilon\}}(x, \omega) d x P(d \omega)
$$

Since the integrand converges to $I_{\{\nabla \xi=0\}}$ when $\varepsilon \rightarrow 0$, then by Lemma (3.1) we have

$$
I_{2}^{\varepsilon} \rightarrow I_{2} \quad \text { as } \quad \varepsilon \rightarrow 0
$$

Let us denote $K_{\varepsilon}(\omega)=\left\{x \mid(x, \omega) \in K_{\varepsilon}\right\}$. Then by the implicit function theorem, for every $\omega$, every $\tau \in \mathbb{R}$ and any $\varepsilon>0$, the set

$$
\Gamma_{\varepsilon}(\tau, \omega)=\left\{x \in K_{\varepsilon}(\omega) \mid \xi(x)=\tau\right\}
$$

is a finite union of $C^{3}$-smooth curves of finite length. We shall denote by $\gamma$ points of a set $\Gamma_{\varepsilon}(\tau, \omega)$ and denote by $d \gamma$ the length-element. Performing the 'co-area change of variables'

$$
K_{\varepsilon}(\omega) \ni x \rightarrow(\tau, \gamma), \quad \tau=\xi(x), \gamma \in \Gamma_{\varepsilon}(\tau, \omega)
$$

we have $d \tau d \gamma=|\nabla \xi| d x_{1} d x_{2}$. So

$$
I_{1}^{\varepsilon}=\mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma_{e}(\tau, \omega)}|\nabla \xi| d \gamma d \tau
$$

and

$$
I_{2}^{\varepsilon}=\mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma_{\varepsilon}(\tau, \omega)}|\nabla \xi|^{-1} d \gamma d \tau
$$

Since the map $\xi$ is $C^{3}$-smooth, then the Sard lemma applies and for every $\omega$ and almost every $\tau \in \mathbb{R}$ the level-set

$$
\begin{equation*}
\Gamma(\tau, \omega)=\left\{x \in \mathbb{T}^{2} \mid \xi(x)=\tau\right\} \tag{3.5}
\end{equation*}
$$

is a $C^{3}$-smooth manifold. So

$$
\begin{equation*}
\int_{\Gamma_{e}(\tau, \omega)}|\nabla \xi| d \gamma \nearrow \int_{\Gamma(\tau, \omega)}|\nabla \xi| d \gamma<\infty \text { as } \varepsilon \rightarrow 0 \tag{3.6}
\end{equation*}
$$

for every $\omega$ and a.e. $\tau$. Hence,

$$
I_{1}^{\varepsilon} \rightarrow \mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi| d \gamma d \tau \text { as } \varepsilon \rightarrow 0
$$

by the monotone convergence theorem. Since also $I_{1}^{\varepsilon} \rightarrow I_{1}$, then

$$
\begin{equation*}
I_{1}=\mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi| d \gamma d \tau \tag{3.7}
\end{equation*}
$$

Similar,

$$
\int_{\Gamma_{e}(\tau, \omega)}|\nabla \xi|^{-1} d \gamma \nearrow \int_{\Gamma(\tau, \omega)}|\nabla \xi|^{-1} d \gamma \leq \infty \text { as } \varepsilon \rightarrow 0
$$

for every $\omega$ and a.e. $\tau$, where we accept the following convention: $0^{-1}=0$. By this convergence,

$$
I_{2}^{\varepsilon} \rightarrow \mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi|^{-1} d \gamma d \tau \text { as } \varepsilon \rightarrow 0
$$

and

$$
I_{2}=\mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi|^{-1} d \gamma d \tau
$$

Now (3.4) implies that

$$
\begin{equation*}
\mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi| d \gamma d \tau=B_{1} \mathbf{E} \int_{\mathbb{R}} g(\tau) \int_{\Gamma(\tau, \omega)}|\nabla \xi|^{-1} d \gamma d \tau \tag{3.8}
\end{equation*}
$$

where both integrals are finite. Since for each $\omega$ the set of critical values $\tau$ of $\xi$ has zero measure, then we can arbitrarily re-define $|\nabla \xi|^{-1}$ in critical points of $\xi$, without changing the integral in the r.h.s.. Below we adopt the following natural convention:

$$
\begin{equation*}
\int_{\Gamma(\tau, \omega)}|\nabla \xi|^{-1} d \gamma=\infty \text { if } \tau \text { is a critical value of } \xi \tag{3.9}
\end{equation*}
$$

(i.e., if $\Gamma(\tau, \omega)$ is not a finite union of smooth curves). Finally, if $\tau$ is a critical value, we define the integral $\int_{\Gamma(\tau, \omega)}|\nabla \xi| d \gamma \leq \infty$ using the limit (3.6). Now the internal integrals in the both sides of (3.8) are defined for all $\omega$ and all $\tau$.

Theorem 3.2. Let $\xi_{v}(t, x)$ be the vorticity of the stationary solution $u_{\nu}(t, x)$.
Then for any $t \geq 0$ and $v>0$ we have

$$
\mathbf{E} \int_{\Gamma(\tau, \omega)}\left|\nabla \xi_{v}\right| d \gamma=B_{1} \mathbf{E} \int_{\Gamma(\tau, \omega)}\left|\nabla \xi_{\nu}\right|^{-1} d \gamma
$$

for a.a. $\tau \in \mathbb{R}$. Here $\Gamma(\tau, \omega)$ is defined in (3.5) (where $\xi(x)=\xi_{v}(t, x)$ ), and we assume (3.9). Moreover, $\int_{\mathbb{R}}\left(\mathbf{E} \int_{\Gamma(\tau, \omega)}\left|\nabla \xi_{\nu}\right| d \gamma\right) d \tau=\frac{1}{2} M_{1}$.

Proof: The first assertion follows from the relation (3.8) which holds for any bounded continuous function $g$. The second assertion follows from (3.7) with $g=1$ since for $g=1$ we have $I_{1}=\frac{1}{2} M_{1}($ see (1.3)).

## 4. BOUNDS FOR EXPONENTIAL MOMENTS

As before, we abbreviate $\xi_{v}(t, x)=\xi(x)$. Choosing in (3.3) $g(v)=v^{2 m}$, where $m$ is a natural number, and denoting $\xi(x)^{m+1}=w(x)$, we get:

$$
\begin{equation*}
\mathbf{E} \int_{\mathbb{T}^{2}}|\nabla w(x)|^{2} d x=B_{1}(m+1)^{2} \mathbf{E} \int_{\mathbb{T}^{2}}|w(x)|^{\frac{2 m}{m+1}} d x \tag{4.1}
\end{equation*}
$$

By the Hölder inequality,

$$
\begin{equation*}
\int|w(x)|^{\frac{2 m}{m+1}} d x \leq(2 \pi)^{\frac{2}{m+1}}\left(\int w(x)^{2} d x\right)^{\frac{m}{m+1}} \tag{4.2}
\end{equation*}
$$

By (3.1), $w(x) \in C^{1}\left(\mathbb{T}^{2}\right)$ for all $\omega$. We wish to estimate the integral of $w(x)^{2}$ in the r.h.s. of (4.2) by its integral Dirichlet. If $m$ was zero, then $w(x)=\xi(x)$ would be a $C^{1}$-smooth function with zero mean-value, and the estimate would follow from the Poincare inequality. Since $m \geq 1$, then what we need is the 'uniform nonlinear Poincaré inequality', given below. A short and nice proof of this result, suggested by M. Struwe and S. Brandle, is given in Appendix.

Lemma 4.1. There exists a constant $C \geq 1$ such that for any $k \in \mathbb{N}$ and any $C^{1}$-smooth function $u(x)$ on $\mathbb{T}^{2}$ with zero mean-value we have

$$
\begin{equation*}
\int u^{2 k} d x \leq C \int\left|\nabla\left(u^{k}\right)\right|^{2} d x \tag{4.3}
\end{equation*}
$$

Combining (4.1)-(4.3) we get:

$$
\mathbf{E} \int|\nabla w(x)|^{2} d x \leq B_{1}(m+1)^{2}(2 \pi)^{\frac{2}{m+1}} C \mathbf{E}\left(\int|\nabla w|^{2} d x\right)^{\frac{m}{m+1}}
$$

Therefore,

$$
\mathbf{E} \int|\nabla w(x)|^{2} d x \leq C_{1}^{m+1}(m+1)^{2(m+1)}
$$

Using (4.3) (with $k=m+1$ ) once again we see that

$$
\begin{equation*}
\mathbf{E} \int \xi(x)^{2(m+1)} d x \leq C_{2}^{m+1}(m+1)^{2(m+1)} \tag{4.4}
\end{equation*}
$$

Since $\xi$ is a homogeneous random field, then the l.h.s. equals $(2 \pi)^{2} \mathbf{E} \xi(x)^{2(m+1)}$. Thus, we have proved that

$$
\begin{equation*}
\left(\mathbf{E}|\xi(x)|^{j}\right)^{1 / j} \leq C j \tag{4.5}
\end{equation*}
$$

for $j=2 r, r=2,3, \ldots$ Since for any random variable $\eta$ the function

$$
(0,1] \ni t \rightarrow \ln \left(\mathbf{E}|\eta|^{1 / t}\right)^{t} \leq \infty
$$

is convex by the Hölder inequality, then (4.5) holds for all $j \in \mathbb{N}$.
For $\sigma>0$ we have $\mathbf{E} e^{\sigma|\xi(x)|}=\sum_{m}\left(\sigma^{m} / m!\right) \mathbf{E}|\xi(x)|^{m}$. As $m!>(m / e)^{m}$ by the Stirling formula, then (4.5) implies that

$$
\mathbf{E} e^{\sigma|\xi(x)|} \leq \sum_{m}(\sigma e C)^{m}
$$

We have got
Theorem 4.2. There exists $\sigma>0$ and $C \geq 1$ such that for any $t \geq 0, x \in \mathbb{T}^{2}$ and $v>0$ we have

$$
\begin{equation*}
\mathbf{E} e^{\sigma\left|\xi_{v}(t, x)\right|} \leq C \tag{4.6}
\end{equation*}
$$

Remark 1. A vector field $u_{v}$ can be recovered from its vorticity $\xi_{v}$ as

$$
\begin{equation*}
u_{v}=D \xi_{v}, \quad D=\left(\frac{\partial}{\partial x_{2}},-\frac{\partial}{\partial x_{1}}\right)^{t}(-\Delta)^{-1} \tag{4.7}
\end{equation*}
$$

where $\Delta$ is the Laplacian, operating on functions on $\mathbb{T}^{2}$ with zero mean-value. The operator $D$ is bounded as a linear map in $L_{2}\left(\mathbb{T}^{2}\right)$ and as an operator in $L_{\infty}\left(T^{2}\right)$. Hence, by the interpolation theorem its norm as an operator in $L_{p}\left(\mathbb{T}^{2}\right)$, $2 \leq p \leq \infty$, is bounded by a $p$-independent constant $C^{\prime}$. Due to this observation, (4.7) and (4.4),

$$
\mathbf{E} \int u_{\nu}(t, x)^{2(m+1)} d x \leq C^{\prime 2(m+1)} C_{2}^{m+1}(m+1)^{2(m+1)}
$$

Therefore

$$
\begin{equation*}
\mathbf{E} e^{\sigma_{1}\left|u_{v}(t, x)\right|} \leq C_{1} \tag{4.8}
\end{equation*}
$$

for some $\sigma_{1}, C_{1}>0$.

Remark 2. Let us abbreviate $u_{\nu}(t, x)=u(x)$. Due to (4.7), $\nabla u=\nabla D \xi$. So $\nabla u(x)$ is obtained from $\xi(x)$ by applying a singular integral operator. Hence, $|\nabla u|_{L_{p}} \leq C_{p}|\xi|_{L_{p}}$ for $1<p<\infty$ by the Calderón-Zydmund theorem. For $2 \leq$ $p<\infty$ the constant $C_{p}$ can be chosen $C_{p}=C p$, e.g., see [17], section II.6. Due to this estimate with $p=2(m+1)$ and (4.4) we have

$$
\mathbf{E} \int|\nabla u(x)|^{2(m+1)} d x \leq C_{3}^{m+1}(m+1)^{4(m+1)}
$$

As before, this inequality implies that $\left(\mathbf{E}|\nabla u(x)|^{j / 2}\right)^{1 / j} \leq C j$ for all $j \in \mathbb{N}$ and all $x$. Therefore,

$$
\begin{equation*}
\mathbf{E} e^{\sigma\left|\nabla u_{v}(t, x)\right|^{1 / 2}} \leq C_{2} \tag{4.9}
\end{equation*}
$$

for any $t, x$, with some $\sigma_{2}, C_{2}>0$.
Remark 3. It is shown in [15] that, along sequences $v_{j} \rightarrow 0$, the processes $u_{v_{j}}(t, x)$ converge in distribution to limiting stationary processes $u_{0}(t, x)$ such that a.a. their trajectories satisfy the Euler equation. Using the Fatou lemma it is easy to check that $u_{0}(t, x)$ and $\xi_{0}(t, x)=\operatorname{curl} u_{0}(t, x)$ satisfy (4.5) and (4.6) (as before, first one has to establish analogies of the inequalities (4.4)). Similar, $u_{0}$ satisfy (4.8) and (4.9).

## 5. APPENDIX

### 5.1. Proof of Lemma (3.1)

Let us consider the space $\mathfrak{A}=[0,1] \times \mathbb{T}^{2} \times \Omega$, given the product sigmaalgebra and the product-measure $\mathcal{P}=L_{t} \times L_{x} \times \mathbf{P}$, where $L_{t}$ is the Lebesgue measure on $[0,1]$ and $L_{x}$ is the normalised Lebesgue measure on $\mathbb{T}^{2}$. Re-defining $\zeta$ and $\xi_{v}$ on a null-set as at the beginning of Section (3) we achieve that (3.1) hold for all $\omega$.

We set $Q=\left\{(t, x, \omega) \in \mathfrak{A} \mid \nabla \xi_{v}=0\right\} \quad$ (as before, $\nabla=\nabla_{x}$ ) and $Q(x)=\{(t, \omega) \mid(t, x, \omega) \in Q\}, \quad Q(t, x)=\{\omega \mid(t, x, \omega) \in Q\}, Q(x, \omega)=$ $\{t \mid(t, x, \omega) \in Q\}$. Since the random function $\xi_{v}(t, x)$ is stationary in $t$ and $x$, then $p:=\mathbf{P}(Q(t, x))$ is independent of $(t, x)$ and $\mathcal{P}(Q)=p$. We have to prove that $p=0$. Assume that, on the contrary, $p>0$.

Let us fix any $x_{0} \in \mathbb{T}^{2}$ and denote by $q^{\omega}$ the set of points of density of $Q\left(x_{0}, \omega\right) \subset[0,1]$. For $t \in[0,1]$ we denote $\pi(t)=\mathbf{P}\left\{t \in q^{\omega}\right\}$. The set $Q^{+}\left(x_{0}\right)=$ $\left\{(t, \omega) \mid \omega \in \Omega, t \in q^{\omega}\right\}$ is measurable as well as the set $Q\left(x_{0}\right)=\{(t, \omega) \mid \omega \in$ $\left.\Omega, t \in Q\left(x_{0}, \omega\right)\right\}$ (since the former may be obtained from the latter in a constructive
way). Therefore we have

$$
\begin{aligned}
\int_{0}^{1} \pi(t) d t & =\int_{0}^{1} \int_{\Omega} I_{Q^{+}\left(x_{0}\right)} d \mathbf{P} d t=\int_{\Omega}\left(\int_{0}^{1} I_{q^{\omega}} d t\right) d \mathbf{P} \\
& =\int_{\Omega}\left(\int_{0}^{1} I_{Q\left(x_{0}, \omega\right)} d t\right) d \mathbf{P}=\left(L_{t} \times \mathbf{P}\right) Q\left(x_{0}\right)=p
\end{aligned}
$$

In particular, there exists $t_{0}<1$ such that $\pi\left(t_{0}\right)>0$.
Lemma 5.1. In $[0,1]$ there exists a converging sequence $t_{n} \searrow t_{0}(n \geq 1)$ such that

$$
\begin{equation*}
\mathbf{P}\left\{\omega \mid\left(t_{n}, x_{0}, \omega\right) \in Q \quad \forall n \geq 0\right\}>0 \tag{5.1}
\end{equation*}
$$

Proof: Since $\pi\left(t_{0}\right)>0$, then $\delta:=\mathbf{P}(S)>0$, where $S=\left\{\omega \mid t_{0} \in q^{\omega}\right\}$. For $0<$ $\tau \leq 1-t_{0}$ we denote $f^{\omega}(\tau)=\tau^{-1} L_{t}\left(q^{\omega} \cap\left[t_{0}, t_{0}+\tau\right]\right) I_{S}(\omega)$. Then

$$
f^{\omega}(\tau) \rightarrow 1 \quad \text { when } \quad \tau \rightarrow 0, \quad \text { for each } \omega \in S
$$

Hence, $\mathbf{E} f^{\omega}(\tau) \rightarrow \delta$ as $\tau \rightarrow 0$, and for any $\varepsilon>0$ there exists $\tau_{\varepsilon}>0$ such that

$$
\mathbf{E} f^{\omega}\left(\tau_{\varepsilon}\right)=\int_{t_{0}}^{t_{0}+\tau_{\varepsilon}} \pi_{S}(\tau) \frac{d \tau}{\tau_{\varepsilon}} \geq \delta(1-\varepsilon), \quad \pi_{S}(\tau)=\mathbf{P}\left\{\tau \in q^{\omega}, \omega \in S\right\}
$$

Accordingly, for any $\varepsilon>0$ we can find $\tau_{\varepsilon}^{\prime} \in\left(0, \tau_{\varepsilon}\right]$ which satisfies $\pi_{S}\left(t_{0}+\tau_{\varepsilon}^{\prime}\right) \geq$ $\delta(1-\varepsilon)$. Now we can use induction to construct a sequence $t_{n} \searrow t_{0}$ with the property $\pi_{S}\left(t_{n}\right) \geq \delta\left(1-2^{-n-1}\right)$. Then $\mathbf{P}\left\{\omega \in S \mid\left(t_{n}, x_{0}, \omega\right) \in Q \quad \forall n\right\} \geq \frac{1}{2} \delta$, and (5.1) follows.

Let us denote

$$
\beta(t)=\nabla \operatorname{curl} \zeta\left(t, x_{0}\right), \quad w(t)=\nabla \xi_{v}\left(t, x_{0}\right)-\sqrt{\nu} \nabla \operatorname{curl} \zeta\left(t, x_{0}\right) .
$$

By (2.2) and (2.3) $\beta(t)$ is a Brownian motion. By (3.1), $\left|w(t)-w\left(t_{0}\right)\right| \leq C^{\omega} \mid t-$ $t_{0} \mid$ for each $\omega$. So

$$
\sqrt{\nu}\left|\beta(t)-\beta\left(t_{0}\right)\right| \leq C^{\omega}\left|t-t_{0}\right|+\left|\nabla \xi_{v}\left(t, x_{0}\right)-\nabla \xi_{v}\left(t_{0}, x_{0}\right)\right| .
$$

This inequality and (5.1) imply that the event

$$
\begin{equation*}
\left\{\left|\beta\left(t_{n}\right)-\beta\left(t_{0}\right)\right| \leq v^{-1 / 2} C^{\omega}\left(t_{n}-t_{0}\right), \quad n=1,2, \ldots\right\} \tag{5.2}
\end{equation*}
$$

has a positive probability. But this is impossible. Indeed, let us consider the event

$$
\begin{equation*}
\bigcap_{N} \bigcup_{n \geq N}\left\{\left|\beta\left(t_{n}\right)-\beta\left(t_{0}\right)\right| \geq \sqrt{t_{n}-t_{0}}\right\} \tag{5.3}
\end{equation*}
$$

We claim that its probability is one. Since (5.2) and (5.3) do not intersect, then the probability of (5.2) must be zero.

It remains to verify the Borel-Cantelli like claim we have made. It suffice to prove that the probability one has the event, obtained from (5.3) by replacing the sequence $\left\{t_{n}\right\}$ by any its subsequence $\left\{t_{n}^{\prime}\right\}$. To do this it suffices to check that

$$
\begin{equation*}
\mathbf{P}\left\{\cap_{n \geq k}\left\{\left|\beta\left(t_{n}^{\prime}\right)-\beta\left(t_{0}\right)\right| \leq \sqrt{t_{n}^{\prime}-t_{0}}\right\}\right\}=0, \tag{5.4}
\end{equation*}
$$

for any $k \geq 1$. The 1.h.s. is bounded by

$$
\begin{equation*}
\mathbf{E} \prod_{n=k}^{K} I_{\left\{\left|\beta\left(t_{n}^{\prime}\right)-\beta\left(t_{0}\right)\right| \leq \sqrt{\left.t_{n}^{\prime}-t_{0}\right\}}\right.}, \tag{5.5}
\end{equation*}
$$

for any fixed $K>k$. If $t_{n+1}-t_{0} \ll t_{n}-t_{0}$, then

$$
\mathbf{E}\left(I_{\left\{\left|\beta\left(t_{n}^{\prime}\right)-\beta\left(t_{0}\right)\right| \leq \sqrt{t_{n}^{\prime}-t_{0}}\right\}} \mid \beta\left(t_{n+1}^{\prime}\right)\right) \leq \mathbf{P}\left\{\left|\beta\left(t_{n}^{\prime}\right)-\beta\left(t_{0}\right)\right| \leq 2 \sqrt{t_{n}^{\prime}-t_{0}}\right\}=: c
$$

for each $\beta\left(t_{n+1}^{\prime}\right)$, satisfying $\left|\beta\left(t_{n+1}^{\prime}\right)-\beta_{t_{0}}\right| \leq \sqrt{t_{n+1}^{\prime}-t_{0}}$. Now, choosing for $\left\{t_{n}^{\prime}\right\}$ a subsequence, converging to $t_{0}$ fast enough and using the Markov property, we see that (5.5) $\leq c^{K-k+1}$. Since $c<1$ and $K$ is any number $>k$, then (5.4) follows.

### 5.2. Proof of Lemma (4.1)

Assume that the lemma's assertion is wrong. Then there exists a sequence of functions $\left\{u_{l}\right\} \subset C^{1}\left(\mathbb{T}^{2}\right)$ with zero mean-value, and a sequence of integers $\left\{p_{l} \geq 1\right\}$ such that for $v_{l}(x)=\operatorname{sgn} u_{l}(x)\left|u_{l}(x)\right|^{p_{l}}$ we have

$$
\begin{equation*}
\int v_{l}^{2} d x \geq l \int\left|\nabla v_{l}\right|^{2} d x \tag{5.6}
\end{equation*}
$$

Without loss of generality we may assume that

$$
\begin{equation*}
\int v_{l}^{2} d x=(2 \pi)^{2} \quad \forall l \tag{5.7}
\end{equation*}
$$

Due to (5.6) and (5.7) the sequence $\left\{v_{l}\right\}$ is bounded in $H^{1}\left(\mathbb{T}^{2}\right)$. So it contains a converging subsequence $\left\{v_{l_{j}}\right\}$ :

$$
\begin{equation*}
v_{l_{j}} \rightarrow v \text { weakly in } H^{1}\left(T^{2}\right) \text { and } v_{l_{j}} \rightarrow v \text { strongly in } L_{2}\left(T^{2}\right) \tag{5.8}
\end{equation*}
$$

By (5.7), $\|v\|_{L_{2}}=2 \pi$. Due to (5.6), $\left\|\nabla v_{l}\right\|_{L_{2}} \leq 2 \pi / \sqrt{l}$. Therefore $\|\nabla v\|_{L_{2}} \leq$ $\liminf \left\|\nabla v_{l_{j}}\right\|_{L_{2}}=0$, and $v(x) \equiv 1$.

Now let us consider the function

$$
\phi(p, u)=\left|\frac{|u|^{p} \operatorname{sgn} u-1}{u-1}\right|, \quad p \geq 1, u \neq 1 .
$$

We claim that

$$
\begin{equation*}
\phi(p, u)>1 / 2 \quad \forall p \geq 1, \forall u \neq 1 \tag{5.9}
\end{equation*}
$$

Indeed, note first that $\phi(1, u) \equiv 1$. It is easy to check that as a function of $p \geq 1$, $\phi(p, u)$ increases with $p$ if $|u|>1$ and if $u \in(0,1)$, so for such $u$ the estimate holds. It remains to consider the case $u \in[-1,0]$. But now

$$
\phi=\frac{|u|^{p}+1}{|u|+1}>\frac{1}{2} .
$$

By (5.9), $\left|u_{l}(x)-1\right| \leq 2\left|v_{l}(x)-1\right|$. Due to this inequality and (5.8) (where $v \equiv 1$ ), we have that $u_{l_{j}} \rightarrow 1$ in $L_{1}$. This contradicts the normalisation $\int u_{l} d x=0$.

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[^1]:    ${ }^{2}$ In fact, this is an equivalent reformulation of the relations.

[^2]:    ${ }^{3}$ See $[1,7,11,12]$ and see [14] for discussion of this result and its development. Recently it was announced in [9] that the stationary measure is unique if (2.3) holds for all $|s| \leq 2$ (previously it was known that the result is true if (2.3) holds for $|s| \leq N$, where $N=N(v)<\infty$ is sufficiently large).

