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Uniform exponential stability of the Ekman spiral

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ABSTRACT. This paper studies stability of the Ekman boundary layer. We utilize a new approach developed by the authors in [12] based on Fourier transformed finite vector Radon measures which yields exponential stability of the Ekman spiral. By this method we can also derive very explicit bounds for solutions of the linearized and the nonlinear Ekman system. For example, we can prove these bounds to be uniform with respect to the angular velocity of rotation which has proved to be relevant for several aspects. Another advantage of this approach is that we obtain well-posedness in classes containing nondecaying vector fields such as almost periodic functions. These outcomes give respect to the nature of boundary layer problems and cannot be obtained by approaches in standard function spaces such as Lebesgue, Bessel-potential, Hölder or Besov spaces.

12010 Mathematics Subject Classification. Primary: 28B05, 28C05, 76D05, Secondary: 76U05, 35Q30,

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1. Introduction and main results

In this note we apply a new approach to rotating boundary layers developed in [12] in order to examine stability of the Ekman boundary layer problem

\[
\begin{aligned}
\partial_t v - \nu \Delta v + \omega e_3 \times v + (v \cdot \nabla)v &= -\nabla q \quad \text{in } (0,T) \times G, \\
\text{div } v &= 0 \quad \text{in } (0,T) \times G, \\
v &= U^E|_{\partial G} \quad \text{on } (0,T) \times \partial G, \\
v|_{t=0} &= v_0 \quad \text{in } G.
\end{aligned}
\]  

(1.1)

Here \( e_3 = (0,0,1)^T \), \( \nu > 0 \) is the viscosity coefficient, and \( \omega \in \mathbb{R} \) is the Coriolis parameter which equals twice the angular velocity of rotation. For \( G \) we will consider simultaneously the half-space \( \mathbb{R}^3_+ \) or a layer, i.e., we have \( G = \mathbb{R}^2 \times D \) with \( D = (0,d) \) and either fixed \( d \in (0,\infty) \) or \( d = \infty \). The vector field \( U^E \) is the so-called Ekman spiral (introduced by the geophysicist V.W. Ekman [8]) given as

\[
U^E(x_3) = U_\infty(1 - e^{-x_3/\delta} \cos(x_3/\delta), e^{-x_3/\delta} \sin(x_3/\delta), 0)^T, \quad x_3 \geq 0.
\]  

(1.2)

System (1.1) is known to be a well-established model for the layer arising in a rotating system (e.g. the earth) between a straight geostrophic flow (e.g. wind) and the surface on which the no slip condition is imposed.

Observe that in the above model rotation about the \( x_3 \)-axis is assumed, whereas \( U_\infty \) denotes the total velocity of the flow, blowing in direction of the \( x_1 \)-axis. The parameter \( \delta \) denotes the layer thickness given by \( \delta = \sqrt{2\nu/|\omega|} \). By geostrophic approximation (see [23]) (1.1) is a reasonable model at least for the upper part of the northern hemisphere. The couple \((U^E, p^E)\) with pressure

\[
p^E(x_2) = -\omega U_\infty x_2
\]

represents a stationary solution of system (1.1). Note that \( U^E(0) = 0 \), i.e. sytem (1.1) is subject to Dirichlet conditions at the lower boundary, and that \( U^E \) is oscillating and nondecaying in tangential direction. We note that remarkable persistent stability of \( U^E \) is observed in geophysical literature.

As for the Ekman problem, the tangentially nondecaying and oscillating behavior is typical for geostrophic boundary layer problems. To give respect to this fact, it seems natural to consider this type boundary layer problems in classes containing nondecaying functions. Hence, the frequently performed \( L^p \) approach for \( 1 < p < \infty \) to the corresponding mathematical models fails in this situation. Giving account to this fact, in [12] an operator theory on spaces of Fourier transformed finite vector-valued Radon measures is developed. These spaces in particular include nondecaying such as almost periodic functions (see Remark 2.9).

A further advantage in dealing with Fourier transformed quantities lies in the fact that all performed calculations and estimations become rather explicit. As a consequence we can derive detailed information on how the solution depends on involved parameters such as time, viscosity, layer thickness, and angular velocity of rotation. In particular, we obtain that the corresponding bounds are uniform in the angular velocity of rotation. This turned out to be relevant for several reasons such as, for instance, the investigation of statistical properties of turbulence, cf. [22, 25]. It also represents the basis for the examination of rapidly oscillating limits as \( \omega \to \infty \),
Mathematically, an approach to stability in time is given in [6] and to asymptotic stability in [15]. These two papers consider the problem in the $L^2$ setting which, of course, does not include nondecaying perturbations of $U^E$. We refer to [14], [4], [5], and [24], [17] for more mathematical literature on the Ekman problem dealing also with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification. For a spectral analysis of the linearized problem we refer to [13], [21]. Local-in-time well-posedness with vanishing Rossby and Ekman numbers and with stratification.
Then the C₀-semigroup \( \exp(-tA_{SCE}) \) on \( FM_{\ell,\sigma}(\mathbb{R}^2, X_2) \) satisfies
\[
(i) \quad \| \exp(-tA_{SCE}) \|_{\mathcal{L}(FM_{\ell,\sigma}(\mathbb{R}^2, X_2))} \leq e^{-2\kappa_\ell t} \quad (t \geq 0),
\]
\[
(ii) \quad \| \nabla \exp(-A_{SCE} - \kappa_\ell \cdot) u_0 \|_{L^2(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2))} \leq \frac{\|u_0\|_{FM_{\ell,\sigma}(\mathbb{R}^2, X_2)}}{\sqrt{\nu(1-\sqrt{2}Re)}},
\]
\[
(iii) \quad \| \nabla \exp(-A_{SCE} - \kappa_\ell \cdot) * f \|_{L^2(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2))} \leq \frac{\|f\|_{L^1(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2))}}{\sqrt{\nu(1-\sqrt{2}Re)}},
\]
for all \( u_0 \in FM_{\ell,\sigma}(\mathbb{R}^2, X_2) \) and \( f \in L^1(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2)) \), where \( g * f(t) := \int_0^t g(t-s)f(s)ds \). In particular, all estimates are uniform in \( \omega \in \mathbb{R} \).

**Remark 1.2.** The physically most relevant case for the layer height is \( d > \delta \).

Based on Theorem 1.1 and a fixed point argument, in Section 4 we derive the following main result for the full nonlinear Ekman problem (1.1).

**Theorem 1.3.** Let the assumptions of Theorem 1.1 be satisfied, and let \( U^E \) be the Ekman spiral given in (1.2). Then for every \( v_0 \in FM_{\ell,\sigma}(\mathbb{R}^2, X_2) + U^E \) such that
\[
\|v_0 - U^E\|_{FM_{\ell,\sigma}(L^2)} < \begin{cases} \pi \nu(1-\sqrt{2}Re)/2^{1/4}3\sqrt{d}, & \text{if } \ell = 0, \ d < \infty, \\ \pi \nu(1-\sqrt{2}Re)\sqrt{7}/2^{1/4}3, & \text{if } \ell > 0, \ d = \infty, \end{cases}
\]
there is a unique global (mild) solution \( v \) of (1.1) satisfying
\[
\exp(2\kappa_\ell \cdot)(v - U^E) \in \text{BC} \left( (0, \infty), FM_{0,\sigma}(\mathbb{R}^2, X_2) \right),
\]
\[
\exp(\kappa_\ell \cdot) \nabla (v - U^E) \in L^2 \left( (0, \infty), FM_{0,\sigma}(\mathbb{R}^2, X_2) \right).
\]
In particular, the Ekman spiral is exponentially stable. More precisely, we have
\[
\|v(t) - U^E\|_{FM_{\ell,\sigma}(\mathbb{R}^2, X_2)} \leq 2 \exp(-2\kappa_\ell t)\|v_0 - U^E\|_{FM_{\ell,\sigma}(\mathbb{R}^2, X_2)} \quad (t \geq 0).
\]
In addition, all estimates above are uniform in \( \omega \in \mathbb{R} \), i.e., with respect to the angular velocity of rotation.
Remark 1.4. By standard bootstrap arguments it can be proved that the solution \( v \) given by Theorem 1.3 enjoys higher regularity in \( FM_\ell (\mathbb{R}^2, X_2) \). By this fact we can recover the pressure via

\[
\nabla q = (I - P)(\nu \Delta v - \omega e_3 \times v - (v \cdot \nabla)v).
\]

Then it can be shown that

\[
(v, q) \in C^\infty((0, \infty) \times \mathbb{R}^2 \times (0, d)),
\]

i.e., \((v, q)\) is the unique classical solution of problem (1.1).

2. Vector-valued Radon measures

Here we introduce notation and recall basic ingredients on \( X \)-valued Radon measures. For basic theory we refer to [7]. The theory related to boundary layers is developed in [12]. There also the proofs of the results listed below can be found.

We use standard notation throughout this article. The symbols \( \mathbb{R}, \mathbb{C}, \mathbb{Z} \) denote the set of reals, complex numbers, and integers, respectively. We also write \( \mathbb{N} = \{1, 2, 3, \ldots\} \) for the naturals and set \( \mathbb{N}_0 := \mathbb{N} \cup \{0\} \). The symbols \( X, Y, Z \) usually denote Banach spaces, whereas \( \mathcal{L}(X, Y) \) stands for the set of bounded linear operators from \( X \) to \( Y \). If \( X = Y \), we write \( \mathcal{L}(X) \). Let \( G \subset \mathbb{R}^n \) be a domain. As usual \( L^2(G, X) \) and \( H^k(G, X) = W^{k,2}(G, X) \) denote \( X \)-valued Lebesgue and Sobolev space respectively. The space \( C^\infty_c(G, X) \) is the set of smooth and compactly supported functions. We will also write \( C(G, X) \), \( BC(G, X) \), and \( BUC(G, X) \) for the space of continuous, the space of bounded and continuous, and the space of bounded and uniformly continuous functions, respectively. The ball in \( \mathbb{R}^n \) centered at \( x_0 \) with radius \( R > 0 \) is denoted by \( B(0, R) \).

The Fourier transformation on the space of rapidly decreasing functions \( S(\mathbb{R}^n) \) in this note is defined as

\[
\hat{u}(\xi) = \mathcal{F}u(\xi) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{-i\xi x} u(x) dx, \quad u \in S(\mathbb{R}^n).
\]

As usual, its extension by duality to the space of tempered distributions \( S'(\mathbb{R}^n, X) := \mathcal{L}'(S(\mathbb{R}^n), X) \) is again denoted by \( \mathcal{F}u \) or \( \hat{u} \) for \( u \in S'(\mathbb{R}^n, X) \).

Next we recall some basic definitions related to \( X \)-valued measures, cf. [7].

Definition 2.1. Let \( X \) be a Banach space, \( \Omega \) be a set, \( \mathcal{A} \) be a \( \sigma \)-algebra over \( \Omega \), and \( \mu : \mathcal{A} \to X \) be a set function.

(i) The function \( \mu \) is called a vector-valued (or \( X \)-valued) measure, if \( \mu(\emptyset) = 0 \) and if it is \( \sigma \)-additive, that is, if it satisfies

\[
\mu\left( \bigcup_{j=1}^\infty A_j \right) = \sum_{j=1}^\infty \mu(A_j)
\]

for all pairwise disjoint sets \( A_j \in \mathcal{A}, j = 1, 2, \ldots \).
(ii) The variation of an $X$-valued measure $\mu$ is defined as

$$|\mu|(O) := \sup \left\{ \sum_{A \in \Pi(O)} \|\mu(A)\|_X : \Pi(O) \subset \mathcal{A} \text{ finite decomposition of } O \right\}.$$ 

for $O \in \mathcal{A}$. (Note that $\Pi(O)$ is a decomposition of $O \in \mathcal{A}$, if $A \cap B = \emptyset$ for all $A, B \in \Pi(O)$ with $A \neq B$ and $\bigcup_{A \in \Pi(O)} A = O$.)

(iii) The quantity $|\mu|(\Omega)$ is called total variation of $\mu$. If $|\mu|(\Omega) < \infty$, then $\mu$ is called finite or of bounded variation.

Next we define $X$-valued Radon measures. For this purpose let $\Omega \subset \mathbb{R}^n$ be open, $\mathcal{A}$ be a $\sigma$-algebra over $\Omega$, and denote by $\mathcal{B}(\Omega)$ the Borel $\sigma$-algebra over $\Omega$. Recall that $\eta : \mathcal{A} \to [0, \infty)$ is a Radon measure if it is Borel regular, that is, if $\mathcal{B}(\Omega) \subseteq \mathcal{A}$ and if for each $A \subseteq \Omega$ there exists a $B \in \mathcal{B}(\Omega)$ such that $A \subseteq B$ and $\eta^*(A) = \eta^*(B)$, where $\eta^*$ denotes the outer measure associated to $\eta$ given by

$$\eta^*(A) := \inf \left\{ \sum_{j=1}^\infty \eta(E_j) : (E_j)_{j \in \mathbb{N}} \subseteq \mathcal{A}, A \subseteq \bigcup_{j=1}^\infty E_j \right\}. \quad (2.1)$$

Also observe that in the sequel we identify a measure $\eta$ by its outer measure, so that $\eta$ is complete in the sense that all subsets $B$ of a set $A \in \mathcal{A}$ satisfying $\eta(A) = 0$ belong to $\mathcal{A}$.

**Definition 2.2.** Let $\Omega \subset \mathbb{R}^n$ be open, $X$ be a Banach space, and $\mathcal{A}$ be a $\sigma$-algebra over $\Omega$. The set function $\mu : \mathcal{A} \to X$ is called a finite $X$-valued Radon measure, if $\mu$ is an $X$-valued measure and if the variation $|\mu|$ is a finite Radon measure. The set of all finite $X$-valued Radon measures is denoted by $M(\Omega, X)$.

From now on assume $X$ to have the Radon-Nikodým property, and $\Omega \subset \mathbb{R}^n$ to be open. By $\rho_\mu \in L^1(\Omega, X, |\mu|)$ we denote the Radon-Nikodým derivative of a measure $\mu \in M(\Omega, X)$ with respect to $(\Omega, \mathcal{A}, |\mu|)$. Then we have the representation

$$\mu(O) = \int_O \rho_\mu \, d|\mu| \quad (O \in \mathcal{A}).$$

Note that by definition each vector Radon measure is well-defined on $\mathcal{B}(\Omega)$. By this fact, for every $\psi \in BC(\Omega, L(Y, Z))$, where $Y$ is another Banach space, its multiplication to an arbitrary $\mu \in M(\Omega, X)$ can be defined as

$$\mu|\psi| := \int_O \psi \rho_\mu \, d|\mu| \quad (O \in \mathcal{A}). \quad (2.2)$$

The properties of this quantity are summarized in

**Lemma 2.3.** ([12, Lemma 2.6]) Let $\Omega \subset \mathbb{R}^n$ be open and let $X, Y, Z$ be Banach spaces having the Radon-Nikodým property. Furthermore, let $\mu \in M(\Omega, X)$ and the functions $\psi, \phi \in BC(\Omega, L(Y, Z))$ be given. Then we have

(i) $|\mu| \psi = |\mu| \|\psi \rho_\mu\|_Y \leq |\mu| \|\psi\|_{L(Y, Z)}$,

(ii) $\mu|\psi| \in M(\Omega, Y)$,

(iii) $(\mu|\psi|) \phi = \mu|\phi \psi|$. 


By the intention to introduce the Fourier transform of vector Radon measures, from now on we assume $\Omega = \mathbb{R}^n$. Note that the fact that the Schwartz space $\mathcal{S}(\mathbb{R}^n)$ of rapidly decreasing functions with its canonical topology is continuously and densely embedded in $\mathcal{L}(C_\infty(\mathbb{R}^n)) := \{v \in C(\mathbb{R}^n) : \lim_{R \to 0} \sup_{B(0,R)} |v(x)| = 0\}$ gives us

$$\mathcal{L}(C_\infty(\mathbb{R}^n), X) \hookrightarrow \mathcal{L}(\mathcal{S}(\mathbb{R}^n), X) = \mathcal{S}'(\mathbb{R}^n, X).$$

Thus, in the sense of the identification

$$\mu \mapsto T_\mu, \quad T_\mu f := \mu \lfloor f(\mathbb{R}^n),$$

we have the embedding

$$\mathcal{M}(\mathbb{R}^n, X) \hookrightarrow \mathcal{S}'(\mathbb{R}^n, X).$$

This allows for the definition of the space of Fourier transformed Radon measures

$$\mathcal{F}\mathcal{M}(\mathbb{R}^n, X) := \{\hat{\mu} : \mu \in \mathcal{M}(\mathbb{R}^n, X)\},$$

which we equip with the canonical norm

$$\|u\|_{\mathcal{F}\mathcal{M}} := \|\mathcal{F}^{-1}u\|_{\mathcal{M}}.$$

Observe that replacing the Fourier transform by its inverse in the definition does not change the value of the norm, i.e., we have $\|\cdot\|_{\mathcal{F}\mathcal{M}} = \|\mathcal{F}\cdot\|_{\mathcal{M}} = \|\mathcal{F}^{-1}\cdot\|_{\mathcal{M}}$. In order to define multipliers with symbols not necessarily continuous at the origin, we also introduce the spaces

$$\mathcal{M}_0(\mathbb{R}^n, X) := \{\mu \in \mathcal{M}(\mathbb{R}^n, X) : \mu(\{0\}) = 0\},$$

that is, the subspace of Radon measures with no point mass at the origin and

$$\mathcal{F}\mathcal{M}_0(\mathbb{R}^n, X) := \{\hat{\mu} : \mu \in \mathcal{M}_0(\mathbb{R}^n, X)\}.$$

Related to exponential stability we introduce further subspaces. These rely on sum-closed frequency sets, which are defined as follows.

**Definition 2.4.** We say that $F \subseteq \mathbb{R}^n$ is a sum-closed frequency set in $\mathbb{R}^n$, if

(i) $F$ is closed,
(ii) $0 \notin F$,
(iii) $F + F := \{x + y : x, y \in F\} \subseteq F \cup \{0\}$.

For a sum-closed frequency set with distance $\ell \sqrt{2} > 0$ from zero in the sequel we write $F_\ell$.\(^1\) For consistency we also set $F_0 = \mathbb{R}^n \setminus \{0\}$ if $\ell = 0$. The class of all sum-closed frequency sets in $\mathbb{R}^n$ is denoted by $\mathcal{F}^n$.

Typical examples of sum-closed frequency sets are (see also [10]):

(i) Countable sum-closed frequency sets in $\mathbb{R}^n$ for which pairwise distances between frequency vectors are uniformly bounded from zero. This case corresponds to almost periodic initial data.

\(^1\)Observe that this differs from the definition of $F_\ell$ in [10] by the factor $\sqrt{2}$. 
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(ii) The set $\mathbb{Z}^n \setminus \{0\}$ or more general

$$F := \left\{ \sum_{j=1}^n m_j a_j ; \ m = (m_1, \ldots, m_n) \in \mathbb{Z}^n \right\} \setminus \{0\},$$

where $a = \{a_1, \ldots, a_n\}$ represents a basis of $\mathbb{R}^n$. This case corresponds to periodic initial data. Indeed, $\supp \hat{u}_0$ is contained in the above $F$ for some $a$ if and only if $u_0$ is periodic. Clearly, this is a special case of (i).

(iii) The set $\{x \in \mathbb{R}^n : x_j \geq \varepsilon\}$ for $j \in \{1, \ldots, n\}$, $\varepsilon > 0$. Note that this example provides non real-valued initial data only.

We also rigorously clarify what we mean by the support of a vector measure.

**Definition 2.5.** Let $X, Y$ be Banach spaces having the Radon-Nikodým property. For $\mu \in M_0(\mathbb{R}^n, X)$ we set

$$N_\mu := \bigcup_{\mathcal{O} \subset \mathbb{R}^n \text{ open, } \mu = 0 \text{ on } \mathcal{O}} \mathcal{O},$$

where we recall that vanishing of a vector-valued measure is canonically defined as

$$\mu = 0 \text{ on } \mathcal{O} \iff \mu(E) = 0 \ (E \subset \mathcal{O}).$$

The support of $\mu$ is defined as

$$\text{supp} \mu := \mathbb{R}^n \setminus N_\mu.$$

**Remark 2.6.** Note that the support of $\mu$ defined above coincides with the support of $\mu$ regarded as a tempered distribution.

In the sequel we will frequently make use of the following observation, mostly without any further notice.

**Remark 2.7.** Let $X, Y$ be Banach spaces having the Radon-Nikodým property. For $\sigma \in \text{BC}(\mathbb{R}^n \setminus \{0\}, \mathcal{L}(X, Y))$ and $\mu \in M_0(\mathbb{R}^n)$ we have

$$\text{supp} \mu | \sigma \subset \text{supp} \mu.$$

**Proof.** Let $\mathcal{O} \subset \mathbb{R}^n$ be open such that $\mu = 0$ on $\mathcal{O}$. Since $M_0(\mathbb{R}^n, X)$ is a Banach space (see Lemma 2.8(i)), the latter is equivalent to $|\mu|(\mathcal{O}) = 0$. Since $\mu|\sigma$ by Lemma 2.3(i) is obviously continuous with respect to $|\mu|$, this yields $\mu|\sigma = 0$ on $\mathcal{O}$. Consequently, we have $N_\mu \subset N_{\mu|\sigma}$ which implies

$$\text{supp} \mu | \sigma = \mathbb{R}^n \setminus N_{\mu|\sigma} \subset \mathbb{R}^n \setminus N_\mu = \text{supp} \mu.$$

\[\square\]

Now let $\ell \geq 0$. For $\ell = 0$ we set $\text{FM}_0(\mathbb{R}^n, X) := \text{FM}_0(\mathbb{R}^n, X)$. For $\ell > 0$ and a sum-closed frequency set $F_\ell \in \mathcal{F}^n$ we also define the space

$$\text{FM}_\ell(\mathbb{R}^n, X) := \{u \in \text{FM}_0(\mathbb{R}^n, X) : \supp \hat{u} \subset F_\ell\}.$$

Next, we list some useful properties of the spaces just introduced.
Lemma 2.8. ([12, Lemma 2.12]) Suppose $X, X_1, X_2$ are Banach spaces having the Radon-Nikodým property and that $X_2 - X_1$ have an embedding constant less or equal to one. Then the following assertions hold.

(i) The spaces $M(R^n, X)$, $M_0(R^n, X)$, and therefore also the spaces $FM(R^n, X)$ and $FM_0(R^n, X)$ for $\ell \geq 0$ are Banach spaces.

(ii) For all $u \in FM(R^n, X_2)$ and $v \in FM(R^n, X_1)$ we have that

$$
\|u \cdot v\|_{FM(R^n, X)} \leq (2\pi)^{-n/2}\|u\|_{FM(R^n, X_2)}\|v\|_{FM(R^n, X_1)},
$$

i.e., $FM(R^n, X_2) \cdot FM(R^n, X_1) \hookrightarrow FM(R^n, X)$. In particular, $(FM(R^n, X), \cdot)$ is an (abelian) algebra (with unit), if $(X, \cdot)$ is an (abelian) algebra (with unit).

(iii) We have

$$
\mathcal{F}L^1(R^n, X) \hookrightarrow FM_0(R^n, X) \hookrightarrow \mathcal{B}_{\infty,1}^0(R^n, X) \hookrightarrow BUC(R^n, X),
$$

where $\mathcal{B}_{\infty,1}^0(R^n, X)$ denotes the homogeneous Besov space.

Remark 2.9. The fact that $\delta_{t_0}a \in M_0(R^n, X)$ for Dirac measures $\delta_{t_0}$, $t_0 \in R^n \setminus \{0\}$, and $a \in X$, gives rise to another interesting class of functions contained in the space $FM_0(R^n, X)$. In fact, every sequence $(a_j)_{j \in \mathbb{N}} \subseteq X$ satisfying $\sum_{j=1}^{\infty} \|a_j\|_X < \infty$ defines for each sequence of frequencies $(\lambda_j)_{j \in \mathbb{N}} \subseteq F\ell$ an element

$$
\left( x \mapsto \sum_{j=1}^{\infty} a_j e^{-i\lambda_j x} \right) \in FM_0(R^n, X),
$$

by the fact that $\sum_{j=1}^{\infty} \delta_{\lambda_j}a_j \in M_0(R^n, X)$. This class of almost periodic functions is significant for applications to rotating boundary layers as explained in the introduction.

For $\sigma \in BC(R^n \setminus \{0\}, \mathcal{L}(X, Y))$ we define

$$
op(\sigma)f := \mathcal{F}^{-1} \hat{f}|\sigma, \quad f \in FM_0(R^n, X).
$$

(2.4)

We recall three results from [12] which allow for a transfer of $L^2$-boundedness to the FM-setting. First, as a consequence of the theory for vector measures developed above we obtain the following multiplier result.

Proposition 2.10. ([12, Proposition 2.13]) Let $\ell \geq 0$. Let $X, Y$ be Banach spaces having the Radon-Nikodým property and suppose that $\sigma \in BC(R^n \setminus \{0\}, \mathcal{L}(X, Y))$. Then $\nop(\sigma)$ as defined in (2.4) is bounded from $FM_0(R^n, X)$ to $FM_0(R^n, Y)$ and we have

$$
\|\nop(\sigma)\|_{\mathcal{L}(FM_0(R^n, X), FM_0(R^n, Y))} = \|\sigma\|_{L^\infty(R^n, \mathcal{L}(X, Y))}.
$$
Remark 2.11. If \( H_1 \) and \( H_2 \) are Hilbert spaces, Plancherel’s theorem implies that the right hand side of the equality in Proposition 2.10 equals the operator norm of \( \text{op}(\sigma) \) in \( \mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(\mathbb{R}^n, H_2)) \). Hence in this case we have

\[
\| \text{op}(\sigma) \|_{\mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(\mathbb{R}^n, H_2))} = \| \sigma \|_{\mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(\mathbb{R}^n, H_2))} = \| \sigma \|_{\mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(\mathbb{R}^n, H_2))}.
\]

In the last part of this section the domain \( \mathbb{R}^n \) is essentially fixed. So we occasionally suppress \( \mathbb{R}^n \) and simply write \( \text{FM}(X) \) instead of \( \text{FM}(\mathbb{R}^n, X) \) and so on. In applications we will often use the fact noted in the above remark for the case that \( H_1 \) and \( H_2 \) are certain \( L^2 \)-spaces. In the same spirit the following lemma will turn out to be helpful.

Lemma 2.12. ([12, Lemma 2.15]) Let \( \ell \geq 0 \), \( J \subset \mathbb{R} \) be an interval, and let \( H_1 \) and \( H_2 \) be Hilbert spaces. Assume that

\[
L \in \mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(J, L^2(\mathbb{R}^n, H_2))) \quad \text{with} \quad \| L \|_{\mathcal{L}(L^2(H_1), L^2(J, L^2(H_2)))} \leq M
\]

is an operator with a symbol \( \sigma_L \) satisfying

\[
\sigma_L \in C\left( \mathbb{R}^n \setminus \{0\}, \mathcal{L}(H_1, L^2(J, H_2)) \right).
\]

Then we have

\[
L \in \mathcal{L}(\text{FM}_\ell(\mathbb{R}^n, H_1), L^2(J, \text{FM}_\ell(\mathbb{R}^n, H_2))), \quad \| L \|_{\mathcal{L}(\text{FM}_\ell(H_1), L^2(J, \text{FM}_\ell(H_2)))} \leq M.
\]

Remark 2.13. In [12] Proposition 2.10 and Lemma 2.12 are established for the case \( \ell = 0 \). However, in view of Remark 2.7, under the action of operators having a suitable symbol representation, \( \text{FM}_\ell(\mathbb{R}^n, X) \) is a closed invariant subspace of \( \text{FM}_0(\mathbb{R}^n, X) \). Therefore these two results readily generalize to the case \( \ell \geq 0 \).

We will also make use of the following general result on operator-valued convolution.

Lemma 2.14. ([12, Lemma 2.16]) Let \( X, Y \) be Banach spaces, \( 1 \leq p \leq \infty \), \( T \in (0, \infty] \), and set \( J = (0, T) \). For \( g \in \mathcal{L}(X, L^p(J, Y)) \) and \( f \in L^1(J, X) \) we have

\[
\left( t \mapsto g \ast f(t) := \int_0^t g(t-s)f(s)ds \right) \in L^p(J, Y)
\]

and

\[
\| g \ast f \|_{L^p(J, Y)} \leq \| g \|_{\mathcal{L}(X, L^p(J, Y))} \| f \|_{L^1(J, X)}. \tag{2.5}
\]

Remark 2.15. Observe that (2.5) is sharper than the usual Young inequality by the fact that \( \| \cdot \|_{\mathcal{L}(X, L^p(J, Y))} \leq \| \cdot \|_{L^p(J, \mathcal{L}(X, Y))} \), but the converse in general is not true. Indeed, (2.5) provides an estimate for singular integral operators, which is not possible with the standard Young’s inequality.

3. Proof of Theorem 1.1

The major advantage of working with \( d < \infty \) or \( \ell > 0 \) lies in the fact that exponential decay of the Stokes-Coriolis-Ekman semigroup can be provided. This follows as a consequence of Poincaré’s inequality, which is valid in these situations.
Let $\ell \geq 0$, $d \in (0, \infty]$ and as before we set $G = \mathbb{R}^2 \times (0, d)$ and $X_2 = L^2((0, d))^3$. Similar to the FM-setting we put

$$L^2_\ell(G)^3 = L^2_\ell(\mathbb{R}^2, X_2) := \{ v \in L^2(\mathbb{R}^2, X_2) : \text{supp} \hat{v} \in F_\ell \}. \quad (1)$$

We emphasize that here $v$ is regarded as a function $v: \mathbb{R}^2 \to X_2$ with $\text{supp} \hat{v} \subset \mathbb{R}^2$.

**Lemma 3.1.** Let either $\ell = 0$, $0 < d < \infty$ or $\ell > 0$, $d = \infty$. Then

$$\|u\|_{L^2_\ell(G)^3} \leq \left\{ \begin{array}{ll} d \sqrt{2} \|\nabla u\|_{L^2_\ell(G)^3}, & \ell = 0, \\ \frac{1}{\ell \sqrt{2}} \|\nabla u\|_{L^2_\ell(G)^3}, & \ell > 0, \end{array} \right.$$ 

for all $u \in L^2_\ell(G)^3$. The assertion remains true, if we replace $L^2_\ell(G)^3$ by the space $FM_\ell(\mathbb{R}^2, X_2)$.

**Proof.** The inequalities in case that $\ell = 0$ (i.e. $d < \infty$) are standard, since we can apply the one-dimensional Poincaré inequality in $x_3$-direction. For $\ell > 0$ we give a proof in the FM-setting. By Plancherel’s theorem the $L^2$-case follows similarly.

We observe that by definition $\text{supp} \hat{u} \subset \mathbb{R}^2 \setminus B(0, \ell \sqrt{2})$. Since $1 \leq |\xi|/\ell \sqrt{2}$ on $\mathbb{R}^2 \setminus B(0, \ell \sqrt{2})$ this gives us in view of Lemma 2.3(i),

$$|\hat{u}|(B(0, R)) = (|\hat{u}|(\{\ell \sqrt{2} \leq |\xi| \leq R\})) \leq \left( |\hat{u}|(\{|\xi|/\ell \sqrt{2}\}) (\{\ell \sqrt{2} \leq |\xi| \leq R\}) \right) = \frac{1}{\ell \sqrt{2}} |\hat{u}|i\xi|(B(0, R)) \quad (R > \ell \sqrt{2}, \ u \in FM_\ell(\mathbb{R}^2, X_2)).$$

Thus we have

$$\|u\|_{FM(X_2)} = \lim_{R \to \infty} |\hat{u}|(B(0, R)) \leq \frac{1}{\ell \sqrt{2}} \lim_{R \to \infty} |\hat{u}|i\xi|(B(0, R)) = \frac{1}{\ell \sqrt{2}} \|\nabla u\|_{FM(X_2)} \quad (u \in FM_\ell(\mathbb{R}^2, X_2)).$$

□

In order to estimate the crucial perturbation arising from the Ekman spiral, also the following Poincaré type inequality will be used, cf. [12, Lemma 3.7]. Its proof is a simple consequence of the fundamental theorem of calculus.

**Lemma 3.2.** We have

$$\|e^{-(\cdot)/\alpha}v\|_{L^2((0, d))} \leq \alpha \left( \int_0^d \frac{e^{-2x} x dx}{\alpha} \right)^{1/2} \|v'\|_{L^2((0, d))} \quad (\alpha, d > 0, \ v \in C^\infty_c((0, d))).$$
Next, we recall existence of the Helmholtz decomposition and the Stokes-Coriolis-Ekman semigroup in the FM-setting. We define solenoidal fields as
\[
\text{FM}_{\ell,\sigma}(\mathbb{R}^{n-1}, X_p) := \left\{ u \in \text{FM}_{\infty}^0 \left( \mathbb{R}^{n-1}, W^{k,p}(D) \right) \mid \text{div } u = 0, \ u^3|_{\partial(\mathbb{R}^{n-1} \times D)} = 0 \right\}
\]
and gradient fields as
\[
G_{\text{FM}} = \{ \nabla p ; p \in L^1_{\text{loc}}(\mathbb{R}^{n-1} \times D), \nabla p \in \text{FM}_{\ell}(\mathbb{R}^{n-1}) \}.
\]
In [12, Lemma 3.4] the following is established.

**Proposition 3.3.** Let \( d \in (0, \infty) \) and \( \ell \geq 0 \). We have the Helmholtz decomposition
\[
\text{FM}_{\ell}(\mathbb{R}^2, X_2) = \text{FM}_{\ell,\sigma}(\mathbb{R}^2, X_2) \oplus G_{\text{FM}}.
\]
The associated Helmholtz projector \( P : \text{FM}_{\ell}(\mathbb{R}^2, X_2) \to \text{FM}_{\ell,\sigma}(\mathbb{R}^2, X_2) \) admits a symbol representation \( \sigma_P = FPF^{-1} \in BC(\mathbb{R}^2 \setminus \{0\}, \mathcal{L}(X_2)) \) such that
\[
\|P\|_{\mathcal{L}(\text{FM}(X_2))} = \|\sigma\|_{L^\infty(\mathbb{R}^2, \mathcal{L}(X_2))} = 1.
\]
Applying \( P \) to the first line of (1.3), the resulting linear operator is given as
\[
A_{\text{SCE}} = A_{\nu} + B_{\omega} + B_E
\]
with
\[
A_{\nu}u = -P\Delta u \quad (\text{Stokes operator}),
\]
\[
B_{\omega}u = \omega Pe_3 \times u \quad (\text{Coriolis part}),
\]
\[
B_Eu = P(U_E \cdot \nabla)u + Pu^3 \partial_3 U_E \quad (\text{Ekman part}).
\]
We call \( A_{\text{SCE}} \) with domain
\[
\mathcal{D}(A_{\text{SCE}}) = \{ u \in \text{FM}_{\ell,\sigma}(\mathbb{R}^{n-1}, X_2) : \partial^\alpha u \in \text{FM}_{\ell}(\mathbb{R}^{n-1}, X_2), \alpha \in \mathbb{N}_{\geq 0}^n, |\alpha| \leq 2, u|_{\partial(\mathbb{R}^{n-1} \times D)} = 0 \}
\]
the Stokes-Coriolis-Ekman operator. In [12, Theorem 3.6] it is proved

**Proposition 3.4.** Let \( d \in (0, \infty] \) and \( \ell \geq 0 \). Then the Stokes-Coriolis-Ekman operator \( A_{\text{SCE}} \) is the generator of a holomorphic \( C_0 \)-semigroup \( (\exp(-tA_{\text{SCE}}))_{t \geq 0} \) on \( \text{FM}_{\ell,\sigma}(\mathbb{R}^2, X_2) \) having a symbol representation
\[
\sigma_{\exp(-t(A_{\text{SCE}} + \lambda_0))) = F\exp(-t(A_{\text{SCE}} + \lambda_0))F^{-1} \in BC(\mathbb{R}^2 \setminus \{0\}, \mathcal{L}(X_2)) \quad (t \geq 0)
\]
for \( \lambda_0 > 0 \) sufficiently large.

**Remark 3.5.** Again we note that, concerning the statements in Propositions 3.3 and 3.4, in [12] only the case \( \ell = 0 \) is treated. However, since we have suitable symbol representations, also here the argument performed in Remark 2.13 applies to the general case.

The space \( L^2_{\ell,\sigma}(G) \) and the Stokes-Coriolis-Ekman operator \( A_{\text{SCE}} \) in \( L^2_{\ell,\sigma}(G) \) are defined accordingly. Now we prove Theorem 1.1 in the \( L^2 \)-setting.
Proposition 3.6. Let \( \nu, U_\infty, d, \ell, \delta, \text{Re}, \kappa_\ell \) satisfy the relations given in Theorem 1.1 and set \( G = \mathbb{R}^2 \times (0, d) \). Then the \( C_0 \)-semigroup \( \exp(-tA_{SCE}) \) on \( L^2_{t,\sigma}(G) \) satisfies

\[
\begin{align*}
(i) \ & \| \exp(-tA_{SCE}) \|_{L^2_{t,\sigma}(G)} \leq e^{-2\kappa_\ell t}, \\
(ii) \ & \| \nabla \exp(-A_{SCE} - \kappa_\ell \cdot t)u \|_{L^2_{t,\sigma}(G)} \leq \frac{1}{\sqrt{\nu(1 - \sqrt{2\text{Re}})}} \| u_0 \|_{L^2_{t,\sigma}(G)},
\end{align*}
\]

for all \( u_0 \in L^2_{t,\sigma}(G) \) and \( t \geq 0 \). In particular, all estimates are uniform in \( \omega \in \mathbb{R} \).

Proof. For \( u_0 \in L^2_{t,\sigma}(G) \), according to Proposition 3.3(b) we may set \( u(t) := e^{-tA_{SCE}}u_0 \). Then \( u \) solves

\[
\begin{align*}
\begin{cases}
u u' + A_{SCE}u & = 0 \quad \text{in } (0, \infty), \\
\hspace{0.5cm} u(0) & = u_0.
\end{cases}
\end{align*}
\]

Multiplying the above equation with \( u \), integrating with respect to \( x \), and taking into account the skew-symmetry of \( B_\omega \) and \( B^1_E \), we obtain

\[
\frac{1}{2} \frac{d}{dt} \| u(t) \|_2^2 + \nu \| \nabla u(t) \|_2^2 + (u^3(t) \partial_3 U^E, u(t)) = 0 \quad (t > 0).
\]

(3.3)

Note that by (1.2) for the derivative of the Ekman spiral we obtain

\[
\partial_3 U^E(x_3) = \frac{U_\infty}{\delta} e^{-x_3/\delta} \begin{pmatrix}
\cos(x_3/\delta) + \sin(x_3/\delta) \\
\cos(x_3/\delta) - \sin(x_3/\delta) \\
0
\end{pmatrix}.
\]

The third term in (3.3) we estimate as follows:

\[
\|(u^3(t) \partial_3 U^E, u(t))\| \leq \sum_{j=1}^2 \| u^3(t) e^{(\cdot)/2\delta}(\partial_3 U^E)^j \|_2 \| e^{(-\cdot)/2\delta} u^j(t) \|_2
\]

\[
\leq \sqrt{2} U_\infty \frac{1}{\delta} \sum_{j=1}^2 \| e^{(-\cdot)/2\delta} u^3(t) \|_2 \| e^{(-\cdot)/2\delta} u^j(t) \|_2
\]

\[
\leq \sqrt{2} U_\infty \delta \| \nabla u(t) \|_2^2,
\]

where we applied twice Lemma 3.2 with \( \alpha = 2\delta \). Inserting this into (3.3) we deduce

\[
\frac{1}{2} \frac{d}{dt} \| u(t) \|_2^2 + \nu \left( 1 - \sqrt{2\text{Re}} \right) \| \nabla u(t) \|_2^2 \leq 0 \quad (t > 0).
\]

(3.4)

Combining estimate (3.4) with Poincaré’s inequality from Lemma 3.1 and keeping in mind that \( \nu - \sqrt{2U_\infty} \delta > 0 \) due to assumption (1.4), we deduce

\[
\frac{d}{dt} \| u(t) \|_2^2 \leq -2\nu \left( 1 - \sqrt{2\text{Re}} \right) \| \nabla u(t) \|_2^2 \leq -4\kappa_\ell \| u(t) \|_2^2 \quad (t > 0).
\]

(3.5)

By virtue of Gronwall’s lemma we therefore obtain

\[
\| u(t) \|_2 \leq e^{-2\kappa_\ell t} \| u_0 \|_2 \quad (t \geq 0).
\]

(3.6)
Thus (i) is proved. Multiplying (3.4) with $e^{2\kappa t}$ and integrating over $t \in \mathbb{R}_+$ yields due to (3.6) that

$$
\int_0^\infty \|e^{\kappa t} \nabla u(s)\|_2^2 ds \leq -\frac{1}{2\nu (1 - \sqrt{2} \text{Re})} \int_0^\infty e^{2\kappa t} \frac{d}{dt} \|u(t)\|_2^2 dt
$$

$$
\leq \frac{1}{2\nu (1 - \sqrt{2} \text{Re})} \left( \|u_0\|_2^2 + 2\kappa \int_0^\infty e^{2\kappa t} \|u(t)\|_2^2 dt \right)
$$

$$
\leq \frac{1}{\nu (1 - \sqrt{2} \text{Re})} \|u_0\|_2^2, \quad (t \geq 0).
$$

Thus the proposition is proved. \(\square\)

We turn to the

**Proof of Theorem 1.1.** The proof is completely analogous to [12, Theorem 7.10]. For the reader’s convenience we repeat the steps. First recall that in [12, Lemma 3.9] it is proved that

$$
\text{Proposition 2.10 then yields (i).}
$$

We turn to the

**Proof of Theorem 1.1.** The proof is completely analogous to [12, Theorem 7.10]. For the reader’s convenience we repeat the steps. First recall that in [12, Lemma 3.9] it is proved that

$$
\text{(1) } (\xi' \mapsto \sigma_{T SCE}(t, \xi')) \in C \left( \mathbb{R}^2 \setminus \{0\}, \mathcal{L}(X_2) \right), \quad t \geq 0, \text{ and}
$$

$$
\text{(2) } (\xi' \mapsto \sigma_{T SCE}(\cdot, \xi')) \in C \left( \mathbb{R}^2 \setminus \{0\}, \mathcal{L}(X_2, L^2((0, T), X_2)) \right)
$$

for every $T \in (0, \infty)$ and where $\sigma_{T SCE} = \mathcal{F} \exp(-tA_{SCE}) \mathcal{F}^{-1}$ denotes the symbol of the Stokes-Coriolis-Ekman semigroup. Note that by Plancherel’s theorem we have

$$
\|\sigma\|_{L^\infty(\mathbb{R}^n, \mathcal{L}(H_1, H_2))} = \|\sigma \rho\|_{\mathcal{L}(L^2(\mathbb{R}^n, H_1), L^2(\mathbb{R}^n, H_1))}
$$

for Hilbert spaces $H_1$, $H_2$. Using this fact, Proposition 3.6(i), Lemma 3.3, and (1) we observe that

$$
\sigma_{T SCE}(t)\sigma_P \in BC \left( \mathbb{R}^2 \setminus \{0\}, \mathcal{L}(X_2) \right)
$$

and

$$
\|\sigma_{T SCE}(t)\sigma_P\|_{L^\infty(\mathbb{R}^2, \mathcal{L}(X_2))} \leq e^{-2\kappa t}, \quad (t \geq 0).
$$

Proposition 2.10 then yields (i).

Next, from Proposition 3.6(ii) we infer

$$
\|\nabla \exp(-(A_{SCE} - \kappa t)\cdot)P\|_{\mathcal{L}(L^2(\mathbb{R}^2, X_2), L^2((0, T), L^2(\mathbb{R}^2, X_2)))} \leq \frac{1}{\sqrt{\nu (1 - \sqrt{2} \text{Re})}}
$$

for all $T > 0$. Setting $H_1 = H_2 = X_2$, $J = (0, T)$, $L = \exp(-(A_{SCE} - \kappa t)\cdot)P$, and $M = 1/\sqrt{\nu (1 - \sqrt{2} \text{Re})}$ we therefore see that relation (ii) is obtained as a consequence of Lemma 2.12 and (2).

In order to see assertion (iii) observe that (ii) implies

$$
\|\nabla \exp(-(A_{SCE} - \kappa t)\cdot)\|
$$

$$
= \|\nabla \exp(-(A_{SCE} - \kappa t)\cdot)\|_{\mathcal{L}(\text{FM}_{\ell,s}(\mathbb{R}^2, X_2), L^2((0, T), \text{FM}_{\ell}(\mathbb{R}^2, X_2)))}
$$

$$
\leq \frac{1}{\sqrt{\nu (1 - \sqrt{2} \text{Re})}} \quad (T > 0).
$$
Thus, with $p = 2$, $g = \nabla \exp(-(A_{SCE} - \kappa \ell) \cdot \cdot)$, $X = FM_{\ell,\sigma}(\mathbb{R}^2, X_2)$, and $Y = FM_{\ell}(\mathbb{R}^2, X_2)$, Lemma 2.14 yields
\[
\| \nabla \exp(-(A_{SCE} - \kappa \ell) \cdot \cdot) * f \|_{L^2((0,T),FM(X_2))} \leq \frac{1}{\sqrt{\nu(1 - \sqrt{2} \Re)} \|f\|_{L^1((0,T),FM(X_2))}} \quad (T > 0)
\]
for all $f \in L^1(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2))$. Hence Theorem 1.1 is proved. \hfill \Box

4. Proof of Theorem 1.3

Again we assume either $\ell = 0$, $0 < d < \infty$ or $\ell > 0$, $d = \infty$. We define the space
\[
\mathbb{E} := \{v \in BC(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2) : \nabla v \in L^2(\mathbb{R}^2, FM_{\ell,\sigma}(\mathbb{R}^2, X_2)),\}
\]
equipped with the norm
\[
\|v\|_\mathbb{E} := \sup_{t > 0} \|e^{2\kappa t}u(t)\|_{FM(X_2)} + \|e^{\kappa t} \cdot \nabla u(t)\|_{L^2(\mathbb{R}^2, FM(X_2))}.
\]
Observe that for $u \in \mathbb{E}$ and $T \in (0, \infty)$ we have $u \in L^2((0, T), FM_{\ell}(\mathbb{R}^2, H^{1}(0, d)^3))$, hence $u(t, x', \cdot)|_{\partial(0, d)}$ and therefore also $u|_{\partial G}$ is well-defined. For fixed initial value $u_0 \in FM_{\ell,\sigma}(\mathbb{R}^2, X_2)$ we further set
\[
B_{u_0} := \{v \in \mathbb{E} : v|_{\partial G} = 0, \ v|_{t=0} = u_0, \ \|v\|_\mathbb{E} \leq M\|u_0\|_{FM(X_2)}\}.
\]
On $B_{u_0}$ we consider the map
\[
Hu(t) := \exp(-tA_{SCE})u_0 - \int_0^t \exp(-(t-s)A_{SCE})P(u(s) \cdot \nabla)u(s)ds, \quad t > 0.
\]
Observe that we have to show that the application of $P$ and of the semigroup to the nonlinear term is well-defined in $FM_{\ell,\sigma}(\mathbb{R}^2, X_2)$. This will be confirmed below after (4.3).

We will show that $H$ is contractive on $B_{u_0}$ for suitable $M > 0$. To this end, we estimate
\[
\|e^{2\kappa t}Hu(t)\|_{FM(X_2)} \leq \|\exp(-(t(A_{SCE} - 2\kappa \ell))u_0\|_{FM(X_2)} + \int_0^\infty \|\exp(-(t-s)(A_{SCE} - 2\kappa \ell))P e^{2\kappa s}u(s) \cdot \nabla u(s)\|_{FM(X_2)}ds \leq \|u_0\|_{FM(X_2)} + \int_0^\infty \|e^{2\kappa s}u(s) \cdot \nabla u(s)\|_{FM(X_2)}ds, \quad (4.2)
\]
where we applied twice Theorem 1.1(i) and then Proposition 3.3. Due to the relation $u(s, x', \cdot)|_{\partial(0, d)} = 0$ we obtain
\[
\|u(s, x', \cdot)\|_{L^\infty(0, d)^3} \leq \sqrt{2} \|u(s, x', \cdot)\|_{L^2(0, d)^3}^{1/2} \|\partial_3 u(s, x', \cdot)\|_{L^2(0, d)^3}^{1/2} \leq \frac{1}{\sqrt{2}} \left(\lambda \|u(s, x', \cdot)\|_2 + \frac{1}{\lambda} \|\partial_3 u(s, x', \cdot)\|_{L^2(0, d)^3}\right)
\]
By the fact that $H^1_0(0,d)^3$, equipped with $\| \cdot \|_\lambda$, and $L^2(0,d)^3$ enjoy the Radon-Nikodým property, we can estimate the nonlinear term thanks to Lemma 2.8(ii) as
\[
\| e^{2\kappa \ell} (u(s) \cdot \nabla) u(s) \|_{FM(X_2)} \\
\leq e^{2\kappa \ell} \frac{3}{2\sqrt{2\pi}} \left( \lambda \| u(s) \|_{FM(L^2)} + \frac{1}{\lambda} \| \nabla u(s) \|_{FM(L^2)} \right) \| \nabla u(s) \|_{FM(X_2)} \\
\leq e^{2\kappa \ell} \frac{3}{2\sqrt{2\pi}} \| u(s) \|^{1/2}_{FM(X_2)} \| \nabla u(s) \|^{3/2}_{FM(X_2)},
\]
where we have set $\lambda = \| \nabla u(s) \|^{1/2}_{FM(X_2)} / \| u(s) \|^{1/2}_{FM(X_2)}$.

In order to be able to further proceed in a unified way for the two different cases $\ell = 0$ and $\ell > 0$ we introduce the following convention:
\[
r := \begin{cases} \\
\sqrt{d}, & \ell = 0, \\
1/\sqrt{\ell}, & \ell > 0.
\end{cases}
\]

By Lemma 3.1 we can then continue the above calculation to the result
\[
\| e^{2\kappa \ell} (u(s) \cdot \nabla) u(s) \|_{FM(X_2)} \leq \frac{3r}{27/4 \pi} \| e^{\kappa s} \nabla u(s) \|_{FM(X_2)}^2. \tag{4.3}
\]

Now we confirm that the nonlinear term even belongs to $FM_\ell(\mathbb{R}^2, X_2)$ such that application of $P$ and afterwards of the Stokes-Coriolis-Ekman semigroup on the space $FM_{\ell,\sigma}(\mathbb{R}^2, X_2)$ is well-defined. For, first we write
\[
(u \cdot \nabla) u = \text{div} \, uu^T.
\]

Since $\text{supp} \hat{u} \subset F_\ell$ and $F_\ell$ is a sum closed frequency set, we have $\text{supp} F uu^T \subset F_\ell \cup \{0\}$, hence that
\[
\text{supp} F \text{div} \, uu^T \subset F_\ell.
\]

(Observe that $F uu^T$ is well-defined as a convolution, cf. [12, Lemma 2.11]; see also [10]). In other words, we have indeed $(u \cdot \nabla) u \in FM_\ell(\mathbb{R}^2, X_2)$.

Plugging (4.3) into (4.2) gives us
\[
\sup_{t > 0} \| e^{2\kappa \ell} H u(t) \|_{FM(X_2)} \leq \| u_0 \|_{FM(X_2)} + \frac{3r}{27/4 \pi} \| e^{\kappa s} \nabla u \|_{L^2(\mathbb{R}^+, FM(X_2))}^2 \\
\leq \| u_0 \|_{FM(X_2)} \left( 1 + M^2 \frac{3r}{27/4 \pi} \| u_0 \|_{FM(X_2)} \right). \tag{4.4}
\]

Next, we derive by utilizing Theorem 1.1(ii) and (iii) that
\[
\| e^{\kappa \ell} \cdot \nabla H u \|_{L^2(\mathbb{R}^+, FM(X_2))} \\
\leq \| \nabla \exp(- (\kappa) u_0) \|_{L^2(\mathbb{R}^+, FM(X_2))}
\]
\[
+ \left( \int_0^\infty \int_0^t \exp(-(t-s)(A_{SCF} - \kappa_\ell)) Pe^{\kappa_\ell s} (u(s) \cdot \nabla) u(s) ds \right) dF_{\alpha}(X_2) dt \right)^{1/2} \\
\leq \frac{1}{\sqrt{\nu(1 - \sqrt{2 Re})}} \left( \| u_0 \|_{FM(X_2)} + \int_0^\infty \| e^{2 \kappa_\ell s} (u(s) \cdot \nabla) u(s) \|_{FM(X_2)} ds \right) \\
\leq \frac{1}{\sqrt{\nu(1 - \sqrt{2 Re})}} \left( \| u_0 \|_{FM(X_2)} + \frac{3r}{2^{7/4} \pi} \| e^{\kappa_\ell s} \nabla u \|_{L^2(R_+, FM(X_2))}^2 \right) \\
\leq \frac{1}{\sqrt{\nu(1 - \sqrt{2 Re})}} \| u_0 \|_{FM(X_2)} \left( 1 + M^2 \frac{3r}{2^{7/4} \pi} \| u_0 \|_{FM(X_2)} \right), \tag{4.5}
\]

where we estimated the nonlinear term again using (4.3). Collecting (4.4) and (4.5) we arrive at

\[
\| Hu \|_E \leq \frac{1}{\sqrt{\nu(1 - \sqrt{2 Re})}} \| u_0 \|_{FM(X_2)} \left( 1 + M^2 \frac{3r}{2^{7/4} \pi} \| u_0 \|_{FM(X_2)} \right)
\]

Thus, choosing

\[
M = \frac{2}{\sqrt{\nu(1 - \sqrt{2 Re})}}
\]

and \( \| u_0 \|_{FM(X_2)} \leq \pi \nu(1 - \sqrt{2 Re})/2^{1/4}3r \) we see that \( H(B_{u_0}) \subseteq B_{u_0} \). Utilizing the expansion

\[
(u \cdot \nabla) u - (v \cdot \nabla)v = ((u - v) \cdot \nabla) u + (v \cdot \nabla)(u - v),
\]

in a very similar way it can be shown that \( H \) is contractive, if the condition on \( u_0 \) is strict. The contraction mapping principle then proves Theorem 1.3.

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**References**