# SPECTRAL PROPERTIES IN $L^{q}$ OF AN OSEEN OPERATOR MODELLING FLUID FLOW PAST A ROTATING BODY 

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#### Abstract

We study the spectrum of a linear Oseen-type operator which arises from equations of motion of a viscous incompressible fluid in the exterior of a rotating compact body. We prove that the essential spectrum consists of an infinite set of overlapping parabolic regions in the left half-plane of the complex plane. The full spectrum coincides with the essential and continuous spectrum if the operator is considered in the whole 3D space. Our approach is based on the Fourier transform in the whole space and the transfer of the results to the exterior domain.


1. Introduction and main results. Suppose that $\mathcal{B}$ is a compact body in $\boldsymbol{R}^{3}$ which is rotating about the $x_{1}$-axis with a constant angular velocity $\omega>0$. Denote by $\Omega(t)$ the exterior of $\mathcal{B}$ at time $t$ and assume that $\Omega(t)$ is a domain with boundary of class $C^{1,1}$. Put $\omega:=\omega \boldsymbol{e}_{1}$, where $\boldsymbol{e}_{1}$ is the unit vector oriented in the direction of the $x_{1}$-axis.

The flow of a viscous incompressible fluid in the exterior of the body $\mathcal{B}$ can be described by the Navier-Stokes system

$$
\begin{aligned}
\partial_{t} \boldsymbol{u}-v \Delta \boldsymbol{u}+(\boldsymbol{u} \cdot \nabla) \boldsymbol{u}+\nabla p & =\boldsymbol{f}, \\
\nabla \cdot \boldsymbol{u} & =0
\end{aligned}
$$

in the space-time region $\left\{(x, t) \in \boldsymbol{R}^{3} \times \mathcal{I} ; t \in \mathcal{I}, x \in \Omega(t)\right\}$, where $\mathcal{I}$ is a time interval. The disadvantage of this description is the variability of the spatial domain $\Omega(t)$. Therefore, many authors use a time-dependent transformation to the Cartesian coordinates, $\boldsymbol{x}^{\prime} \equiv\left(x_{1}^{\prime}, x_{2}^{\prime}, x_{3}^{\prime}\right)$, connected with the rotating body and defined by $\boldsymbol{x}^{\prime}=O(\omega t) \boldsymbol{x}$, where

$$
O(\phi)=\left(\begin{array}{ccc}
1 & 0 & 0  \tag{1.1}\\
0 & \cos \phi & \sin \phi \\
0 & -\sin \phi & \cos \phi
\end{array}\right)
$$

If we also put

$$
\begin{aligned}
& \boldsymbol{u}(\boldsymbol{x}, t)=O^{T}(\omega t) \boldsymbol{u}^{\prime}\left(\boldsymbol{x}^{\prime}, t\right)=O^{T}(\omega t) \boldsymbol{u}^{\prime}(O(\omega t) \boldsymbol{x}, t), \\
& p(\boldsymbol{x}, t)=p^{\prime}\left(\boldsymbol{x}^{\prime}, t\right)=p^{\prime}(O(\omega t) \boldsymbol{x}, t), \\
& \boldsymbol{f}(\boldsymbol{x}, t)=O^{T}(\omega t) \boldsymbol{f}^{\prime}\left(\boldsymbol{x}^{\prime}, t\right)=O^{T}(\omega t) \boldsymbol{f}^{\prime}(O(\omega t) \boldsymbol{x}, t),
\end{aligned}
$$

then the functions $\boldsymbol{u}^{\prime}$ and $p^{\prime}$ satisfy the system

$$
\begin{align*}
\partial_{t} \boldsymbol{u}^{\prime}-v \Delta^{\prime} \boldsymbol{u}^{\prime}-\left(\boldsymbol{\omega} \times \boldsymbol{x}^{\prime}\right) \cdot \nabla^{\prime} \boldsymbol{u}^{\prime}+\boldsymbol{\omega} \times \boldsymbol{u}^{\prime}+\left(\boldsymbol{u}^{\prime} \cdot \nabla^{\prime}\right) \boldsymbol{u}^{\prime}+\nabla^{\prime} p^{\prime} & =\boldsymbol{f}^{\prime}, \\
\nabla^{\prime} \cdot \boldsymbol{u}^{\prime} & =0 \tag{1.2}
\end{align*}
$$

in $\Omega(0) \times \mathcal{I}$, where $\Omega(0)$ is $t$-independent. (Here, $\Delta^{\prime}$ and $\nabla^{\prime}$ denote the operators, acting in the variable $\boldsymbol{x}^{\prime}$.) Among a series of results on qualitative properties of the system (1.2) and related linear problems, let us mention Hishida [21], [22], [23], Galdi [16], [17], Farwig, Hishida, Müller [8], Farwig [4], [5], Farwig, Hishida [6], [7], Galdi, Kyed [18], [19], Geissert, Heck, Hieber [20], Farwig, Neustupa [13], [14], Farwig, Nečasová, Neustupa [11], [12], Farwig, Krbec, Nečasová [9], [10].

If $\boldsymbol{u}(\boldsymbol{x}, t)$ tends to the constant velocity $\gamma \boldsymbol{e}_{1}$ for $|\boldsymbol{x}| \rightarrow \infty$, then $\boldsymbol{u}^{\prime}\left(\boldsymbol{x}^{\prime}, t\right) \rightarrow \gamma \boldsymbol{e}_{1}$ for $\left|\boldsymbol{x}^{\prime}\right| \rightarrow \infty$ because the transformation does not affect the vector $\boldsymbol{e}_{1}$. In this case, it is advantageous to write $\boldsymbol{u}^{\prime}=\boldsymbol{v}^{\prime}+\gamma \boldsymbol{e}_{1}$ and to deal with a new system for the unknown function $v^{\prime}$ :

$$
\begin{align*}
\partial_{t} \boldsymbol{v}^{\prime}-v \Delta^{\prime} \boldsymbol{v}^{\prime}-\left(\boldsymbol{\omega} \times \boldsymbol{x}^{\prime}\right) \cdot \nabla^{\prime} \boldsymbol{v}^{\prime}+\boldsymbol{\omega} \times \boldsymbol{v}^{\prime}+\gamma \partial_{1}^{\prime} \boldsymbol{v}^{\prime}+\left(\boldsymbol{v}^{\prime} \cdot \nabla^{\prime}\right) \boldsymbol{v}^{\prime}+\nabla^{\prime} p^{\prime} & =\boldsymbol{f}^{\prime} \\
\nabla^{\prime} \cdot \boldsymbol{v}^{\prime} & =0 \tag{1.3}
\end{align*}
$$

in $\Omega(0) \times \mathcal{I}$.
In order to have a simple notation, we will further denote $\Omega:=\Omega(0)$. Moreover, since we will work only in the coordinates $\boldsymbol{x}^{\prime}$, we omit the primes and we denote the new coordinate system by $\boldsymbol{x}$ instead of $\boldsymbol{x}^{\prime}$. By analogy, we also omit the primes in the notation of functions.

Similarly, as the properties of the steady linear Stokes system (respectively the so called Stokes operator) play an important role in the qualitative theory of the Navier-Stokes equations, the qualitative analysis of the system (1.3) is based on properties of the steady linear Oseen-type problem

$$
\begin{align*}
\nu \Delta \boldsymbol{v}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}-\nabla p & =\boldsymbol{f} & & \text { in } \Omega, \\
\nabla \cdot \boldsymbol{v} & =0 & & \text { in } \Omega,  \tag{1.4}\\
\boldsymbol{v} & =\mathbf{0} & & \text { on } \partial \Omega .
\end{align*}
$$

This problem can be written in the form of one operator equation $A_{\gamma}^{\omega} \boldsymbol{v}=\boldsymbol{f}$, where $A_{\gamma}^{\omega}$ is the Oseen-type operator

$$
\begin{equation*}
A_{\gamma}^{\omega} \boldsymbol{v}=P_{q} v \Delta \boldsymbol{v}+P_{q}\left[(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}\right] \tag{1.5}
\end{equation*}
$$

with the domain

$$
\begin{equation*}
D\left(A_{\gamma}^{\omega}\right):=\left\{\boldsymbol{v} \in W^{2, q}(\Omega)^{3} \cap W_{0}^{1, q}(\Omega)^{3} \cap L_{\sigma}^{q}(\Omega) ;(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v} \in L^{q}(\Omega)^{3}\right\} \tag{1.6}
\end{equation*}
$$

in the function space $L_{\sigma}^{q}(\Omega), 1<q<\infty$; this function space as well as the Helmholtz projection $P_{q}$ in $L^{q}(\Omega)^{3}$ will be defined below. We shall further treat the domain $D\left(A_{\gamma}^{\omega}\right)$ as a Banach space with the norm

$$
\begin{equation*}
\|\boldsymbol{v}\|_{D\left(A_{\gamma}^{\omega}\right)}:=\|\boldsymbol{v}\|_{2, q}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}\|_{q} . \tag{1.7}
\end{equation*}
$$

Let us note that especially the presence of the term $(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}$ changes substantially the character of the operator $A_{\gamma}^{\omega}$, in comparison with the Stokes operator $A_{0}^{0}$ or with the Oseen operator $A_{\gamma}^{0}$. For instance, while the operators $A_{0}^{0}$ and $A_{\gamma}^{0}$ generate analytic semigroups in the space $L_{\sigma}^{q}(\Omega)$, the operator $A_{\gamma}^{\omega}$ generates only a $C_{0}$-semigroup (see Hishida [22] when $\gamma=0$, $q=2$, Geissert, Heck, Hieber [20] when $\gamma=0,1<q<\infty$, and Shibata [28] when $\gamma \neq 0$, $1<q<\infty)$.

The information on the spectrum of the linear operator $A_{\gamma}^{\omega}$ plays a fundamental role in studies of (1.4) and (1.3). The cases $q=2$ and $\gamma=0$ or $\gamma \neq 0$ were treated in our papers [13], [14], and the case $1<q<\infty, \gamma=0$ was studied in our paper [11], [12], (together with Š. Nečasová). In this paper, we consider the case $1<q<\infty$ and $\gamma \neq 0$ and we study the spectrum of the operator $A_{\gamma}^{\omega}$. The main difference, in comparison with [13], [14], is that $L_{\sigma}^{q}(\Omega)$ (with general $\left.1<q(\neq 2)<\infty\right)$ is not a Hilbert space and the Fourier expansion with respect to the angular variable (in an appropriate cylindrical coordinate system, whose axis is the $x_{1}$-axis), which was the main tool used in [13], [14], cannot be effectively applied now. Thus, our approach is based on the Fourier transform in the whole space and on a transfer of the results to the exterior domain. We have already used the same idea in the papers [11], [12] in the case $\gamma=0$. The case $\gamma \neq 0$, discussed in this paper, however, brings new effects which strongly influence the shape of the spectrum. Without loss of generality, we assume that $\gamma>0$.

Concerning recent results in this field, we mention the papers [9], [10] where weighted $L^{q}$-estimates for stationary solutions of the linear whole space problem in the Stokes as well as Oseen case are derived. Shibata [28] proved that the Oseen operator $A_{\gamma}^{\omega}$ generates a $C^{0}$ semigroup satisfying the usual $L^{q}-L^{r}$-estimates. Moreover, he found detailed estimates of the "localized" resolvent near and on the imaginary axis. Similar results for the Stokes case are known from Hishida, Shibata [24].

We shall use the usual function spaces and norms:

- The norm in $L^{q}(\Omega)^{3}$ is denoted by $\|\cdot\|_{q}=\|\cdot\|_{q ; \Omega}$.
- $W_{0}^{1, q}(\Omega)$ is the subspace of the Sobolev space $W^{1, q}(\Omega)$ consisting of functions vanishing on $\partial \Omega$ in the sense of traces.
- The norm in $W^{k, q}(\Omega)^{3}, k \in N$, is denoted by $\|\cdot\|_{k, q}=\|\cdot\|_{k, q ; \Omega}$.
- $C_{0, \sigma}^{\infty}(\Omega)$ denotes the linear space of all divergence-free vector fields from $C_{0}^{\infty}(\Omega)^{3}$.
- $L_{\sigma}^{q}(\Omega)$ is the closure of $C_{0, \sigma}^{\infty}(\Omega)$ in $L^{q}(\Omega)^{3}$ and coincides with the space of all divergence-free (in the sense of distributions) vector fields $\boldsymbol{u} \in L^{q}(\Omega)^{3}$ such that $\boldsymbol{u} \cdot \boldsymbol{n}=0$ on $\partial \Omega$ in the sense of traces ([15], pp. 111-115); here $\boldsymbol{n}$ is the outer normal vector on $\partial \Omega$.
- $P_{q}$ denotes the projection of $L^{q}(\Omega)^{3}$ onto $L_{\sigma}^{q}(\Omega)$, associated with the Helmholtz decomposition

$$
L^{q}(\Omega)^{3}=L_{\sigma}^{q}(\Omega) \oplus\left\{\nabla p \in L^{q}(\Omega)^{3} ; p \in W_{\mathrm{loc}}^{1, q}(\Omega)\right\}
$$

Let us recall definitions and basic properties from spectral theory of linear operators. Assume that $X$ is a Banach space with the norm $\|\cdot\|, X^{*}$ is its dual and $T$ is a closed linear
operator in $X$. We denote by $D(T)$ the domain and by $R(T)$ the range of $T$ and we assume that $D(T)$ is dense in $X$. This guarantees that the adjoint operator $T^{*}$ exists.

- $\operatorname{nul}(T)$ and $\operatorname{def}(T)$ denote the nullity and the deficiency of the operator $T$, respectively. If $R(T)$ is closed then $\operatorname{nul}(T)=\operatorname{def}\left(T^{*}\right)$ and $\operatorname{def}(T)=\operatorname{nul}\left(T^{*}\right)$ (see e.g. Kato [25, p. 234]).
- $\operatorname{nul}^{\prime}(T)$ and $\operatorname{def}^{\prime}(T):=\operatorname{nul}^{\prime}\left(T^{*}\right)$ denote the approximate nullity and the approximate deficiency of $T$, respectively. We recall that $\operatorname{nul}^{\prime}(T)$ is the maximum integer $m(m=+\infty$ being permitted) with the property that to each $\varepsilon>0$ there exists an $m$-dimensional linear manifold $M_{\varepsilon}$ in $D(T)$ such that $\|T \boldsymbol{v}\|<\varepsilon$ for all $\boldsymbol{v} \in M_{\varepsilon},\|\boldsymbol{v}\|=1$. Note that $\operatorname{nul}(T) \leq$ $\operatorname{nul}^{\prime}(T)$ and $\operatorname{def}(T) \leq \operatorname{def}^{\prime}(T)$ and the equalities hold if $R(T)$ is closed. On the other hand, if $R(T)$ is not closed then $\operatorname{nul}^{\prime}(T)=\operatorname{def}^{\prime}(T)=\infty$. The identity $\operatorname{nul}^{\prime}(T)=\infty$ is equivalent to the existence of a non-compact sequence $\left\{\boldsymbol{u}_{n}\right\}$ on the unit sphere in $X$ such that $T \boldsymbol{u}_{n} \rightarrow \mathbf{0}$ for $n \rightarrow \infty$ (see [25, p. 233]).
- We say that $T$ is a Fredholm operator if its range $R(T)$ is closed in $X$ and both the numbers nul $(T)$ and $\operatorname{def}(T)$ are finite.
- The operator $T$ is called semi-Fredholm if the range $R(T)$ is closed in $X$ and at least one of the numbers nul $(T)$ and $\operatorname{def}(T)$ is finite. Consequently, $T$ is semi-Fredholm if and only if at least one of the numbers nul' $(T)$ and $\operatorname{def}^{\prime}(T)$ is finite.
- We denote by $\rho(T)$ the resolvent set of $T$, by $\sigma(T)=\boldsymbol{C} \backslash \rho(T)$ the spectrum of $T$ and by $\sigma_{\text {ess }}(T)$ the essential spectrum of $T$. Recall that $\rho(T)$ is the set of all $\lambda \in \boldsymbol{C}$ such that $R(T-\lambda I)=X$ and the operator $T-\lambda I$ has a bounded inverse in $X$. Thus, $\operatorname{nul}(T-\lambda I)=\operatorname{nul}^{\prime}(T-\lambda I)=\operatorname{def}(T-\lambda I)=\operatorname{def}^{\prime}(T-\lambda I)=0$ for $\lambda \in \rho(T)$. Moreover, $\sigma_{\text {ess }}(T)=\{\lambda \in \boldsymbol{C} ; T-\lambda I$ is not semi-Fredholm $\}$. Both $\sigma(T)$ and $\sigma_{\text {ess }}(T)$ are closed subsets of $\boldsymbol{C}$ and $\sigma_{\text {ess }}(T) \subset \sigma(T)$.
- Let us also mention that $\sigma(T)=\sigma_{\mathrm{p}}(T) \cup \sigma_{\mathrm{c}}(T) \cup \sigma_{\mathrm{r}}(T)$, where the sets $\sigma_{\mathrm{p}}(T), \sigma_{\mathrm{c}}(T)$ and $\sigma_{\mathrm{r}}(T)$ are called the point spectrum, the continuous spectrum and the residual spectrum of $T$, respectively. They are mutually disjoint and they are defined in this way:
(a) $\sigma_{\mathrm{p}}(T):=\{\lambda \in C ; \operatorname{nul}(T-\lambda I)>0\}$.
(b) $\sigma_{\mathrm{c}}(T)$ is the set of $\lambda \in \boldsymbol{C}$ such that $\operatorname{nul}(T-\lambda I)=0, R(T-\lambda I)$ is dense in $X$, but $R(T-\lambda I) \neq X$. In this case, $R(T-\lambda I)$ is not closed in $X$, which implies that $\operatorname{def}(T-\lambda I)=\operatorname{def}^{\prime}(T-\lambda I)=\operatorname{nul}^{\prime}(T-\lambda I)=\infty$.
(c) $\sigma_{\mathrm{r}}(T)$ is the set of $\lambda \in \boldsymbol{C}$ such that nul $(T-\lambda I)=0$ and the range $R(T-\lambda I)$ is not dense in $X$.

Obviously, $\sigma_{\mathrm{c}}(T) \subset \sigma_{\text {ess }}(T)$. There are no generally valid relations between $\sigma_{\mathrm{p}}(T), \sigma_{\mathrm{r}}(T)$ on one hand and $\sigma_{\text {ess }}(T)$ on the other hand. However, any point on the boundary of $\sigma(T)$ belongs to $\sigma_{\text {ess }}(T)$ unless it is an isolated point of $\sigma(T)$ (see [25, p. 244]).

- The so-called approximate point spectrum $\sigma_{\mathrm{ap}}(T)$ of $T$ consists of all points $\lambda \in \boldsymbol{C}$ such that there exists a sequence $\left\{\boldsymbol{u}_{n}\right\}$ in $D(T)$ such that $\left\|\boldsymbol{u}_{n}\right\|=1$ and $(T-\lambda I) \boldsymbol{u}_{n} \rightarrow \mathbf{0}$ as $n \rightarrow \infty$. Obviously, $\operatorname{nul}^{\prime}(T-\lambda I)>0$ for $\lambda \in \sigma_{\mathrm{ap}}(T)$, which implies that $\lambda \in \sigma(T)$, and $\sigma_{\mathrm{p}}(T) \subset \sigma_{\mathrm{ap}}(T)$. Finally, if $\lambda \in \sigma_{\mathrm{c}}(T)$ then $\operatorname{nul}^{\prime}(T-\lambda I)=\infty$, which also implies that there exists a sequence $\left\{\boldsymbol{u}_{n}\right\}$ with the properties required in the definition of $\sigma_{\mathrm{ap}}(T)$. Hence
$\lambda \in \sigma_{\mathrm{ap}}(T)$. We have thus shown that $\sigma_{\mathrm{p}}(T) \cup \sigma_{\mathrm{c}}(T) \subset \sigma_{\mathrm{ap}}(T) \subset \sigma(T)$ and $\sigma_{\mathrm{ap}}(T) \cup$ $\sigma_{\mathrm{r}}(T)=\sigma(T)$. We note that the approximate spectrum has been introduced for the sake of completeness, but will not be used in this paper.

The main theorems of this paper concern the concrete operator $A_{\gamma}^{\omega}$ in the space $L_{\sigma}^{q}(\Omega)$ :
Theorem 1.1. Let $1<q<\infty$ and $\Omega=\boldsymbol{R}^{3}$. Then
(i) $\sigma\left(A_{\gamma}^{\omega}\right)=\sigma_{\mathrm{c}}\left(A_{\gamma}^{\omega}\right)=\sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)=\Lambda_{\gamma}^{\omega}$, where

$$
\Lambda_{\gamma}^{\omega}:=\left\{\lambda=\alpha+\mathrm{i} \beta+\mathrm{i} k \omega \in \boldsymbol{C} ; \alpha, \beta \in \boldsymbol{R}, k \in \boldsymbol{Z}, \alpha \leq-\nu \beta^{2} / \gamma^{2}\right\} .
$$

(ii) If $q=2$ then $A_{\gamma}^{\omega}$ is a normal operator in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)\left(=L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)\right)$.

THEOREM 1.2. Let $1<q<\infty$ and $\Omega \subset \boldsymbol{R}^{3}$ be an exterior domain with boundary of class $C^{1,1}$. Then the spectrum of $A_{\gamma}^{\omega}$ lies in the left complex half plane $\{\lambda \in \boldsymbol{C} ; \operatorname{Re} \lambda \leq$ $0\}$ and consists of the essential spectrum $\sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)=\Lambda_{\gamma}^{\omega}$ and possibly a set $\Gamma$ of isolated eigenvalues $\lambda \in \boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$ with $\operatorname{Re} \lambda<0$ and finite algebraic multiplicity, which can cluster only at points of $\sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)$. The set $\Gamma$ of such isolated eigenvalues is independent of $q \in(1, \infty)$.

The set $\Lambda_{\gamma}^{\omega}$ is a union of infinitely many equally shifted filled parabolas in the left halfplane of $\boldsymbol{C}$ (see Figure 1). Theorem 1.1 is proved in Section 3. The proof of Theorem 1.2 is given in Section 4.

The question whether the identities of Theorem 1.1 (i) also hold in the case when $\Omega$ is an exterior domain in $\boldsymbol{R}^{3}$ is open. The reason consists in the application of the Fourier transform, which is a useful tool in $\boldsymbol{R}^{3}$ but cannot be used in exterior domains. Moreover, it is open whether the set $\Gamma$ of isolated eigenvalues is empty or non-empty.


Figure 1. The shape of set $\Lambda_{\gamma}^{\omega}$ in the complex domain $\boldsymbol{C}$.
2. Preliminary results. The domain $D\left(A_{\gamma}^{\omega}\right)$ of the operator $A_{\gamma}^{\omega}$ is dense in $L_{\sigma}^{q}(\Omega)$, because $C_{0, \sigma}^{\infty}(\Omega) \subset D\left(A_{\gamma}^{\omega}\right)$ and $C_{0, \sigma}^{\infty}(\Omega)$ is dense in $L_{\sigma}^{q}(\Omega)$. Hence the adjoint operator $\left(A_{\gamma}^{\omega}\right)^{*}$ exists as a linear operator in $L_{\sigma}^{q^{\prime}}(\Omega)$, where $q^{\prime}=q /(q-1)$. The next lemma brings more information on both the operators $A_{\gamma}^{\omega}$ and $\left(A_{\gamma}^{\omega}\right)^{*}$.

Lemma 2.1. The operator $A_{\gamma}^{\omega}$ is closed in $L_{\sigma}^{q}(\Omega)$ and generates a $C_{0}$-semigroup in $L_{\sigma}^{q}(\Omega)$. Its adjoint operator is

$$
\begin{equation*}
\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{v}=P_{q^{\prime}} \Delta \boldsymbol{v}+P_{q^{\prime}}\left[-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}+\boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v}\right] \tag{2.1}
\end{equation*}
$$

with domain

$$
\begin{equation*}
D\left(\left(A_{\gamma}^{\omega}\right)^{*}\right)=\left\{\boldsymbol{v} \in W^{2, q^{\prime}}(\Omega)^{3} \cap W_{0}^{1, q^{\prime}}(\Omega)^{3} \cap L_{\sigma}^{q^{\prime}}(\Omega) ;(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v} \in L^{q^{\prime}}(\Omega)^{3}\right\} \tag{2.2}
\end{equation*}
$$

It is a closed operator in $L_{\sigma}^{q^{\prime}}(\Omega)$ and generates a $C_{0}$-semigroup in $L_{\sigma}^{q^{\prime}}(\Omega)$.
PROOF. The fact that $A_{\gamma}^{\omega}$ is a generator of a $C_{0}$-semigroup in $L_{\sigma}^{q}(\Omega)$ follows from [28, Theorem 1.1]. It also implies that $A_{\gamma}^{\omega}$ is a closed operator in $L_{\sigma}^{q}(\Omega)$ and that $R\left(A_{\gamma}^{\omega}-\zeta I\right)=$ $L_{\sigma}^{q}(\Omega)$ for all $\zeta>0$ sufficiently large. Let us denote by $T_{\gamma}{ }^{\omega}$ the operator on the right hand side of (2.1) with the domain given by (2.2):

$$
D\left(T_{\gamma}^{\omega}\right)=\left\{\boldsymbol{v} \in W^{2, q^{\prime}}(\Omega)^{3} \cap W_{0}^{1, q^{\prime}}(\Omega)^{3} \cap L_{\sigma}^{q^{\prime}}(\Omega) ;(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v} \in L^{q^{\prime}}(\Omega)^{3}\right\}
$$

By analogy with $A_{\gamma}^{\omega}$, the operator $T_{\gamma}^{\omega}$ is closed in $L_{\sigma}^{q^{\prime}}(\Omega)$ and $R\left(T_{\gamma}^{\omega}-\zeta I\right)=L_{\sigma}^{q^{\prime}}(\Omega)$ if $\zeta>0$ is sufficiently large. It is easy to verify that the operators $A_{\gamma}^{\omega}$ and $T_{\gamma}^{\omega}$ are adjoint to each other in the sense of Kato [25, p. 167]; hence $T_{\gamma}^{\omega} \subset\left(A_{\gamma}^{\omega}\right)^{*}$. In order to show that $T_{\gamma}^{\omega}=\left(A_{\gamma}^{\omega}\right)^{*}$, we need to verify that $T_{\gamma}^{\omega}$ is the maximal operator adjoint to $A_{\gamma}^{\omega}$. Suppose that $\boldsymbol{v} \in D\left(\left(A_{\gamma}^{\omega}\right)^{*}\right)$ and put $\boldsymbol{f}=\left(\zeta I-\left(A_{\gamma}^{\omega}\right)^{*}\right) \boldsymbol{v}$. Since $\boldsymbol{f} \in R\left(T_{\gamma}^{\omega}-\zeta I\right)$, there exists $\boldsymbol{w} \in D\left(T_{\gamma}^{\omega}\right)$ such that $f=\left(T_{\gamma}^{\omega}-\zeta I\right) w$. Hence $\left(\left(A_{\gamma}^{\omega}\right)^{*}-\zeta I\right) \boldsymbol{v}=\left(T_{\gamma}^{\omega}-\zeta I\right) w$. Multiplying both sides of this identity by $\boldsymbol{u} \in D\left(A_{\gamma}^{\omega}\right)$ and integrating in $\Omega$, we arrive at

$$
\int_{\Omega} \boldsymbol{v} \cdot\left(A_{\gamma}^{\omega}-\zeta I\right) \boldsymbol{u} \mathrm{d} \boldsymbol{x}=\int_{\Omega} \boldsymbol{w} \cdot\left(A_{\gamma}^{\omega}-\zeta I\right) \boldsymbol{u} \mathrm{d} \boldsymbol{x} .
$$

Since this holds for all $\boldsymbol{u} \in D\left(A_{\gamma}^{\omega}\right)$, we get $\boldsymbol{v}=\boldsymbol{w} \in D\left(T_{\gamma}^{\omega}\right)$; thus $\left(A_{\gamma}^{\omega}\right)^{*}=T_{\gamma}^{\omega}$. As for $A_{\gamma}^{\omega}$, we conclude that $\left(A_{\gamma}^{\omega}\right)^{*}$ generates a $C_{0}$-semigroup in $L_{\sigma}^{q^{\prime}}(\Omega)$ and is closed.

Lemma 2.2. There exist constants $c_{1}>0$ and $c_{2}>0$ such that, if $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and $f \in L_{\sigma}^{q}(\Omega)$ satisfy the equation $A_{\gamma}^{\omega} v=f$, then

$$
\begin{equation*}
\|\boldsymbol{v}\|_{2, q}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}\|_{q} \leq c_{1}\|\boldsymbol{f}\|_{q}+c_{2}\|\boldsymbol{v}\|_{q} \tag{2.3}
\end{equation*}
$$

Proof. If $\Omega=\boldsymbol{R}^{3}$ then (2.3) follows from [4, Theorem 1.1] and an interpolation argument. Now consider an exterior domain $\Omega \subset \boldsymbol{R}^{3}$ of class $C^{1,1}$. Let $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and $f \in L_{\sigma}^{q}(\Omega)$ satisfy the equation $A_{\gamma}^{\omega} v=f$. Then there exists a pressure function $p$ such that

$$
\nabla p=v \Delta \boldsymbol{v}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}-\boldsymbol{f} \in L^{q}(\Omega)^{3} .
$$

For simplicity, we assume that in the following all pressure functions have a vanishing mean on $\Omega_{R}=\Omega \cap B_{R}(\mathbf{0})$, i.e., $\int_{\Omega_{R}} p \mathrm{~d} \boldsymbol{x}=0$, where $R>0$ is chosen such that $\boldsymbol{R}^{3} \backslash B_{R-1}(\mathbf{0}) \subset$ $\Omega$. Let $\eta \in C_{0}^{\infty}\left(\boldsymbol{R}^{3}\right)$ denote a cut-off function with values in $[0,1]$ such that

$$
\eta(\boldsymbol{x})= \begin{cases}0 & \text { for }|\boldsymbol{x}| \leq R-1, \\ 1 & \text { for }|\boldsymbol{x}| \geq R\end{cases}
$$

Then $(\eta \boldsymbol{v}, \eta p)$ can be considered as a solution of the whole space problem

$$
\nu \Delta \boldsymbol{u}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{u}-\boldsymbol{\omega} \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}-\nabla \tilde{p}=\boldsymbol{f}_{1}, \quad \nabla \cdot \boldsymbol{u}=g_{1} \quad \text { in } \boldsymbol{R}^{3},
$$

where

$$
\boldsymbol{f}_{1}=\eta \boldsymbol{f}+2 v \nabla \eta \cdot \nabla \boldsymbol{v}+v \boldsymbol{v} \Delta \eta+((\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \eta) \boldsymbol{v}-\gamma\left(\partial_{1} \eta\right) \boldsymbol{v}-p \nabla \eta, \quad g_{1}=\boldsymbol{v} \cdot \nabla \eta .
$$

Note that $\boldsymbol{f}_{1}$ coincides with $\eta \boldsymbol{f}$ up to perturbation terms with support in $\operatorname{supp}(\nabla \eta) \subset \bar{\Omega}_{R}$. Similarly, supp $g_{1} \subset \bar{\Omega}_{R}$. By [4, Theorem 1.1] there exists a solution ( $\left.\boldsymbol{u}, \tilde{p}\right)$ satisfying the estimate

$$
\left\|\nabla^{2} \boldsymbol{u}\right\|_{q}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{u}-\boldsymbol{\omega} \times \boldsymbol{u}\|_{q}+\|\nabla \tilde{p}\|_{q} \leq c\left(\|\boldsymbol{f}\|_{q}+\|\boldsymbol{v}\|_{1, q ; \Omega_{R}}+\|p\|_{q ; \Omega_{R}}\right) .
$$

Moreover, by the uniqueness assertion in [4, Theorem 1.1], ( $\eta \boldsymbol{v}, \eta p$ ) satisfies the same estimate so that we get the inequality

$$
\begin{align*}
& \left\|\nabla^{2}(\eta \boldsymbol{v})\right\|_{q}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla(\eta \boldsymbol{v})\|_{q}+\|\nabla(\eta p)\|_{q} \\
& \quad \leq c\left(\|\boldsymbol{f}\|_{q}+\|\boldsymbol{v}\|_{q}+\|\boldsymbol{v}\|_{1, q ; \Omega_{R}}+\|p\|_{q ; \Omega_{R}} .\right. \tag{2.4}
\end{align*}
$$

Next we consider $((1-\eta) \boldsymbol{v},(1-\eta) p)$ as a solution of the Stokes problem

$$
\nu \Delta \boldsymbol{u}-\nabla \tilde{p}=\boldsymbol{f}_{2}, \quad \operatorname{div} \boldsymbol{u}=g_{2} \quad \text { in } \Omega_{R}, \quad \boldsymbol{u}=\mathbf{0} \quad \text { on } \partial \Omega_{R},
$$

where $g_{2}=-\boldsymbol{v} \cdot \nabla \eta$ and

$$
\boldsymbol{f}_{2}=(1-\eta)\left[\boldsymbol{f}-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}+\boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v}\right]-2 v \nabla \eta \cdot \nabla \boldsymbol{v}-v \boldsymbol{v} \Delta \eta+p \nabla \eta .
$$

Hence classical a priori estimates for the Stokes system in the bounded domain $\Omega_{R}$ and the precise form of $\boldsymbol{f}_{2}$ imply that

$$
\begin{equation*}
\left\|\nabla^{2}((1-\eta) \boldsymbol{v})\right\|_{q ; \Omega_{R}}+\|\nabla(1-\eta) p\|_{q ; \Omega_{R}} \leq c\left(\|\boldsymbol{f}\|_{q}+\|\boldsymbol{v}\|_{1, q ; \Omega_{R}}+\|p\|_{q ; \Omega_{R}}\right) . \tag{2.5}
\end{equation*}
$$

Summing (2.4) and (2.5) and using an interpolation estimate for $\nabla \boldsymbol{v}$, we obtain the estimate

$$
\begin{equation*}
\left\|\nabla^{2} \boldsymbol{v}\right\|_{q}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}\|_{q}+\|\nabla p\|_{q} \leq c\left(\|\boldsymbol{f}\|_{q}+\|\boldsymbol{v}\|_{q}+\|p\|_{q ; \Omega_{R}}\right) \tag{2.6}
\end{equation*}
$$

with a constant $c>0$ independent of $\boldsymbol{f}, \boldsymbol{v}, p$.
Now we will prove (2.3) by contradiction, using (2.6). Suppose that, for every $n \in N$, there exist $\boldsymbol{v}_{n} \in D\left(A_{\gamma}^{\omega}\right)$ and $p_{n}$ with $\nabla p_{n} \in L^{q}(\Omega)^{3}$ such that

$$
\begin{equation*}
1=\left\|\nabla^{2} \boldsymbol{v}_{n}\right\|_{q}+\left\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}_{n}\right\|_{q}+\left\|\nabla p_{n}\right\|_{q} \geq n\left(\left\|\boldsymbol{f}_{n}\right\|_{q}+\left\|\boldsymbol{v}_{n}\right\|_{q}\right) \tag{2.7}
\end{equation*}
$$

where $\boldsymbol{f}_{n}=A_{\gamma}^{\omega} \boldsymbol{v}_{n}$. Then we find a subsequence of ( $\boldsymbol{v}_{n}, p_{n}$ ), again denoted by ( $\boldsymbol{v}_{n}, p_{n}$ ), and $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and $p$ with $\nabla p \in L^{q}(\Omega)^{3}$ such that in the weak sense

$$
\begin{equation*}
\nabla^{2} \boldsymbol{v}_{n} \rightharpoonup \nabla^{2} \boldsymbol{v}, \quad(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}_{n} \rightharpoonup(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}, \quad \nabla p_{n} \rightharpoonup \nabla p \quad \text { in } L^{q}(\Omega), \tag{2.8}
\end{equation*}
$$

and in the strong sense

$$
\begin{equation*}
f_{n} \rightarrow \mathbf{0}, \quad \boldsymbol{v}_{n} \rightarrow \mathbf{0} \quad \text { in } L^{q}(\Omega), \quad p_{n} \rightarrow p \quad \text { in } L^{q}\left(\Omega_{R}\right), \quad \int_{\Omega_{R}} p \mathrm{~d} \boldsymbol{x}=0 \tag{2.9}
\end{equation*}
$$

as $n \rightarrow \infty$, and

$$
v \Delta \boldsymbol{v}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}-\nabla p=0, \quad \operatorname{div} \boldsymbol{v}=0 \quad \text { in } \Omega, \quad \boldsymbol{v}=\mathbf{0} \quad \text { on } \partial \Omega .
$$

By (2.8), (2.9) we conclude that $\boldsymbol{v}=0$ and $\nabla p=\mathbf{0}$ in $\Omega$. Since the limit function $p$ was chosen with $\int_{\Omega_{R}} p \mathrm{~d} \boldsymbol{x}=0$, we have $p=0$ in $\Omega$. Using now (2.6) with $\boldsymbol{v}_{n}, p_{n}, \boldsymbol{f}_{n}$ instead of $\boldsymbol{v}, p, \boldsymbol{f}$, we observe that the left-hand side equals one, while the right-hand side tends to zero as $n \rightarrow \infty$. This is the contradiction to the assumption that (2.3) was wrong.

Let us further denote by $D_{0}\left(A_{\gamma}^{\omega}\right)$ the subspace of $D\left(A_{\gamma}^{\omega}\right)$ which contains only functions that have a compact support in $\bar{\Omega}$.

LEMMA 2.3. $D_{0}\left(A_{\gamma}^{\omega}\right)$ is a core of the operator $A_{\gamma}^{\omega}$, i.e., the graph of the restriction of the operator $A_{\gamma}^{\omega}$ to $D_{0}\left(A_{\gamma}^{\omega}\right)$ is dense in the graph of $A_{\gamma}^{\omega}$ in the norm of $L_{\sigma}^{q}(\Omega) \times L_{\sigma}^{q}(\Omega)$.

Proof. Let $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and $\boldsymbol{f}=A_{\gamma}^{\omega} \boldsymbol{v}$. We will show that $\left[\boldsymbol{v}, A_{\gamma}^{\omega} \boldsymbol{v}\right]$ can be approximated by a sequence of elements $\left[\boldsymbol{v}^{n}, A_{\gamma}^{\omega} \boldsymbol{v}^{n}\right.$ ] where $\boldsymbol{v}^{n} \in D_{0}\left(A_{\gamma}^{\omega}\right), n \in \boldsymbol{N}$.

Let $\eta \in C_{0}^{\infty}\left(\boldsymbol{R}^{3}\right)$ be radially symmetric, with values in $[0,1]$, such that

$$
\eta(\boldsymbol{x})= \begin{cases}1 & \text { for }|\boldsymbol{x}| \leq 1 \\ 0 & \text { for } 2 \leq|\boldsymbol{x}|\end{cases}
$$

Denote $K_{1}:=\left\{\boldsymbol{x} \in \boldsymbol{R}^{3} ; 1<|\boldsymbol{x}|<2\right\}$ and, more generally, for $r>0$, let $K_{r}:=\{\boldsymbol{x} \in$ $\left.\boldsymbol{R}^{3} ; r<|\boldsymbol{x}|<2 r\right\}$. Due to Bogovskij [2], there exists a bounded linear operator $\mathfrak{B}$ : $L^{q}\left(K_{1}\right) \rightarrow W_{0}^{1, q}\left(K_{1}\right)^{3}$ such that div $\mathfrak{B} f=f$ for all $f \in L^{q}\left(K_{1}\right)$ satisfying $\int_{K_{1}} f \mathrm{~d} \boldsymbol{x}=0$. The operator $\mathfrak{B}$ is bounded from $W_{0}^{1, q}\left(K_{1}\right)$ to $W_{0}^{2, q}\left(K_{1}\right)^{3}$ as well.

Let $n \in \boldsymbol{N}$ be so large that $\boldsymbol{R}^{3} \backslash B_{n}(\mathbf{0}) \subset \Omega$. Put $\boldsymbol{v}^{n}(\boldsymbol{x}):=\eta(\boldsymbol{x} / n) \boldsymbol{v}(\boldsymbol{x})-\boldsymbol{V}^{n}(\boldsymbol{x})$ with the correction term $\boldsymbol{V}^{n}(\boldsymbol{x})$ being equal to $\boldsymbol{U}^{n}(\boldsymbol{x} / n)$, where

$$
\boldsymbol{U}^{n}(\boldsymbol{y})= \begin{cases}\mathfrak{B}[\nabla \eta(\boldsymbol{y}) \cdot \boldsymbol{v}(n \boldsymbol{y})] & \text { for } \boldsymbol{y} \in K_{1}, \\ \mathbf{0} & \text { for } \boldsymbol{y} \in \boldsymbol{R}^{3} \backslash K_{1}\end{cases}
$$

The function $\boldsymbol{v}^{n}$ is divergence-free, it coincides with $\boldsymbol{v}$ in $\Omega \cap B_{n}(\mathbf{0})$ and its support is a subset of the closure of $\Omega \cap B_{2 n}(\mathbf{0})$. Due to the continuity of the operator $\mathfrak{B}$ from $L^{q}\left(K_{1}\right)$ to $W_{0}^{1, q}\left(K_{1}\right)^{3}$ and from $W_{0}^{1, q}\left(K_{1}\right)$ to $W_{0}^{2, q}\left(K_{1}\right)^{3}$, the function $\boldsymbol{U}^{n}$ satisfies the estimates

$$
\begin{aligned}
\left\|\nabla \boldsymbol{U}^{n}\right\|_{q ; K_{1}} & \leq C\|\nabla \eta \cdot \boldsymbol{v}(n \cdot)\|_{q ; K_{1}} \leq C\|\boldsymbol{v}(n \cdot)\|_{q ; K_{1}}=C n^{-3 / q}\|\boldsymbol{v}\|_{q ; K_{n}}, \\
\left\|\nabla^{2} \boldsymbol{U}^{n}\right\|_{q ; K_{1}} & \leq C\|\nabla(\nabla \eta \cdot \boldsymbol{v}(n \cdot))\|_{q ; K_{1}} \leq C\|\boldsymbol{v}(n \cdot)\|_{q ; K_{1}}+C\|\nabla \boldsymbol{v}(n \cdot)\|_{q ; K_{1}} \\
& =C n^{-3 / q}\|\boldsymbol{v}\|_{q ; K_{n}}+C n^{1-3 / q}\|\nabla \boldsymbol{v}\|_{q ; K_{n}} .
\end{aligned}
$$

This means that

$$
\left\|\nabla \boldsymbol{V}^{n}\right\|_{q ; K_{n}}=n^{-1+3 / q}\left\|\nabla \boldsymbol{U}^{n}\right\|_{q ; K_{1}} \leq \frac{C}{n}\|\boldsymbol{v}\|_{q ; K_{n}}
$$

$$
\left\|\nabla^{2} \boldsymbol{V}^{n}\right\|_{q ; K_{n}}=n^{-2+3 / q}\left\|\nabla^{2} \boldsymbol{U}^{n}\right\|_{q ; K_{1}} \leq \frac{C}{n^{2}}\|\boldsymbol{v}\|_{q ; K_{n}}+\frac{C}{n}\|\nabla \boldsymbol{v}\|_{q ; K_{n}}
$$

Using also the fact that $\boldsymbol{V}^{n}(\boldsymbol{x})=\mathbf{0}$ for $|\boldsymbol{x}|=n$, we derive that

$$
\left\|\boldsymbol{V}^{n}\right\|_{q ; K_{n}} \leq C n\left\|\nabla \boldsymbol{V}^{n}\right\|_{q ; K_{n}} \leq C\|\boldsymbol{v}\|_{q ; K_{n}}
$$

All generic constants $C$ are independent of $n$. Thus, $\left\|\boldsymbol{V}^{n}\right\|_{2, q} \rightarrow 0$ as $n \rightarrow \infty$. The same result also holds on the $\|\cdot\|_{2, q}$-norm of the difference $\eta(\boldsymbol{x} / n) \boldsymbol{v}(\boldsymbol{x})-\boldsymbol{v}(\boldsymbol{x})$. Consequently,

$$
\begin{equation*}
\boldsymbol{v}^{n} \rightarrow \boldsymbol{v} \quad \text { in } W^{2, q}(\Omega)^{3} \text { for } n \rightarrow \infty \tag{2.10}
\end{equation*}
$$

Furthermore, since $(\omega \times \boldsymbol{x}) \cdot \nabla \eta(\boldsymbol{x} / n)=0$, we have

$$
\begin{align*}
& \left\|A_{\gamma}^{\omega} \boldsymbol{v}-A_{\gamma}^{\omega} \boldsymbol{v}^{n}\right\|_{q ; \Omega} \\
& \leq C\left(\left\|\boldsymbol{v}-\boldsymbol{v}_{n}\right\|_{2, q ; \Omega}+\left\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}^{n}\right\|_{q ; \Omega}\right) \\
& \leq C\left(\left\|\boldsymbol{v}-\boldsymbol{v}_{n}\right\|_{2, q ; \Omega}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot(1-\eta(\boldsymbol{x} / n)) \nabla \boldsymbol{v}\|_{q ; \Omega}+\left\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{V}^{n}\right\|_{q ; \Omega}\right)  \tag{2.11}\\
& \leq C\left(\left\|\boldsymbol{v}-\boldsymbol{v}_{n}\right\|_{2, q ; \Omega}+\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}\|_{q ; \Omega \backslash B_{n}(0)}+n\left\|\nabla \boldsymbol{V}^{n}\right\|_{q ; \Omega}\right) \\
& \rightarrow 0 \text { for } n \rightarrow \infty
\end{align*}
$$

We can now observe from (2.10) and (2.11) that

$$
\left[\boldsymbol{v}^{n}, A_{\gamma}^{\omega} \boldsymbol{v}^{n}\right] \rightarrow\left[\boldsymbol{v}, A_{\gamma}^{\omega} \boldsymbol{v}\right] \quad \text { in } L_{\sigma}^{q}(\Omega) \times L_{\sigma}^{q}(\Omega)
$$

as $n \rightarrow \infty$. The proof is completed.
3. The case $\Omega=\boldsymbol{R}^{3}$. If $\Omega=\boldsymbol{R}^{3}, 1<q<\infty$ and $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$, then the terms $v \Delta \boldsymbol{v}$, $(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}$ and $\gamma \partial_{1} \boldsymbol{v}$ belong to $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. Hence the projection $P_{q}$ in (1.5) can be omitted so that

$$
\begin{equation*}
A_{\gamma}^{\omega} v=v \Delta \boldsymbol{v}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v} \tag{3.1}
\end{equation*}
$$

By analogy, the adjoint operator, as an operator in $L_{\sigma}^{q^{\prime}}\left(\boldsymbol{R}^{3}\right)$, can be simplified to

$$
\begin{equation*}
\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{v}=v \Delta \boldsymbol{v}-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}+\boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v} \tag{3.2}
\end{equation*}
$$

The next lemma provides an information on solutions of the resolvent equation

$$
\begin{equation*}
A_{\gamma}^{\omega} \boldsymbol{v}-\lambda \boldsymbol{v}=\boldsymbol{f} \tag{3.3}
\end{equation*}
$$

for $\boldsymbol{f} \in L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. Recall the definition of the set $\Lambda_{\gamma}^{\omega}$ from Theorem 1.1.
LEMMA 3.1. Suppose that $\lambda \in \boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$. There exists a constant $c_{3}=c_{3}(\lambda, q)>0$ such that if $\boldsymbol{f} \in L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$ and $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ satisfy the resolvent equation (3.3) then

$$
\begin{equation*}
\|\boldsymbol{v}\|_{q} \leq c_{3}\|\boldsymbol{f}\|_{q} \tag{3.4}
\end{equation*}
$$

Proof. The linear space $D_{0}\left(A_{\gamma}^{\omega}\right)$ is a core of $A_{\gamma}^{\omega}-\lambda I$ due to Lemma 2.3. Thus, it is sufficient to prove (3.4) only for $v \in D_{0}\left(A_{\gamma}^{\omega}\right)$.

Equation (3.3) can be written in the form

$$
\begin{equation*}
v \Delta v+(\omega \times \boldsymbol{x}) \cdot \nabla v-\omega \times v-\gamma \partial_{1} v-\lambda v=f \tag{3.5}
\end{equation*}
$$

Due to the geometry of the problem, it is reasonable to use the cylindrical coordinates $\left(x_{1}, r, \varphi\right)$ in $\boldsymbol{R}^{3}$ with the axis being the $x_{1}$-axis; then $r^{2}=x_{2}^{2}+x_{3}^{2}$. The term $(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}$ in (3.5), which equals $\omega\left(-x_{3} \partial_{2} \boldsymbol{v}+x_{2} \partial_{3} \boldsymbol{v}\right)$, can be simplified to

$$
\begin{equation*}
(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}=\omega \partial_{\varphi} \boldsymbol{v} \tag{3.6}
\end{equation*}
$$

We shall denote by $\mathcal{F}$ the Fourier transform, by $\mathcal{F}_{-1}$ its inverse, by $\widehat{ }$ Fourier images of functions, by $\boldsymbol{\xi}=\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$ their Cartesian Fourier variables, and we put $s=|\boldsymbol{\xi}|$. Applying $\mathcal{F}$ to (3.5), we obtain

$$
\begin{equation*}
\left(\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}\right) \widehat{\boldsymbol{v}}-\omega \partial_{\phi} \widehat{\boldsymbol{v}}+\omega \times \widehat{\boldsymbol{v}}=-\widehat{\boldsymbol{f}} \tag{3.7}
\end{equation*}
$$

Here $\partial_{\phi} \widehat{\boldsymbol{v}}$ denotes the angular derivative

$$
\partial_{\phi} \widehat{\boldsymbol{v}}=\left(\boldsymbol{e}_{1} \times \boldsymbol{\xi}\right) \cdot \nabla \widehat{\boldsymbol{v}} \equiv-\xi_{3} \frac{\partial \widehat{\boldsymbol{v}}}{\partial \xi_{2}}+\xi_{2} \frac{\partial \widehat{\boldsymbol{v}}}{\partial \xi_{3}}
$$

when using cylindrical coordinates $\left(\xi_{1}, \rho, \phi\right)$ in the space of the Fourier variables. The equation $\operatorname{div} \boldsymbol{v}=0$ (following from the fact that $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ ) leads to the condition $\mathbf{i} \boldsymbol{\xi} \cdot \widehat{\boldsymbol{v}}=0$. Now $\widehat{\boldsymbol{v}}$ can be considered to be a solution of the first order ordinary differential equation (3.7) with respect to the angular variable $\phi$. Writing $\widehat{\boldsymbol{v}}$ in the form

$$
\widehat{\boldsymbol{v}}\left(\rho, \phi, \xi_{1}\right)=O(-\phi) \widehat{\boldsymbol{w}}\left(\rho, \phi, \xi_{1}\right)
$$

where $O$ is the matrix from (1.1), one verifies that

$$
\omega \partial_{\phi} \widehat{\boldsymbol{v}}=O(-\phi) \omega \partial_{\phi} \widehat{\boldsymbol{w}}+\boldsymbol{\omega} \times[O(-\phi) \widehat{\boldsymbol{w}}]
$$

Hence (3.7) is equivalent to the equation

$$
\begin{equation*}
-\omega \partial_{\phi} \widehat{\boldsymbol{w}}+\left(\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}\right) \widehat{\boldsymbol{w}}=-O(\phi) \widehat{\boldsymbol{f}} \tag{3.8}
\end{equation*}
$$

or using the definition

$$
a(\xi)=\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}, \quad s=|\boldsymbol{\xi}|
$$

to the inhomogeneous ordinary first order linear differential equation with respect to $\phi$

$$
\omega \partial_{\phi} \widehat{\boldsymbol{w}}-a(\boldsymbol{\xi}) \widehat{\boldsymbol{w}}=O(\phi) \widehat{\boldsymbol{f}}
$$

Its solution $\widehat{\boldsymbol{w}}$ satisfies

$$
\begin{aligned}
& \widehat{\boldsymbol{w}}\left(\xi_{1}, \rho, \phi+2 \pi\right) \\
& \quad=\mathrm{e}^{2 \pi a(\xi) / \omega} \widehat{\boldsymbol{w}}\left(\xi_{1}, \rho, \phi\right)+\frac{1}{\omega} \int_{0}^{2 \pi} \mathrm{e}^{(2 \pi-t) a(\xi) / \omega} O(t+\phi) \widehat{\boldsymbol{f}}\left(\xi_{1}, \rho, t+\phi\right) \mathrm{d} t
\end{aligned}
$$

Since $\widehat{\boldsymbol{w}}$ is $2 \pi$-periodic in the variable $\phi$, we have

$$
\begin{equation*}
\widehat{\boldsymbol{w}}\left(\xi_{1}, \rho, \phi\right)=\frac{1}{\omega} \int_{0}^{2 \pi} \frac{\mathrm{e}^{(2 \pi-t) a(\xi) / \omega}}{1-\mathrm{e}^{2 \pi a(\xi) / \omega}} O(t+\phi) \widehat{\boldsymbol{f}}\left(\xi_{1}, \rho, t+\phi\right) \mathrm{d} t \tag{3.9}
\end{equation*}
$$

and consequently

$$
\begin{equation*}
\widehat{\boldsymbol{v}}\left(\xi_{1}, \rho, \phi\right)=\frac{1}{\omega} \int_{0}^{2 \pi} \frac{\mathrm{e}^{-t a(\xi) / \omega}}{\mathrm{e}^{-2 \pi a(\xi) / \omega}-1} O(t) \widehat{\boldsymbol{f}}\left(\xi_{1}, \rho, t+\phi\right) \mathrm{d} t \tag{3.10}
\end{equation*}
$$

Returning to the Cartesian variables $\boldsymbol{\xi}=\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$ in (3.10), we obtain

$$
\begin{align*}
\widehat{\boldsymbol{v}}(\boldsymbol{\xi}) & =\frac{1}{\omega} \int_{0}^{2 \pi} \Psi(\lambda, \boldsymbol{\xi}, t) O(t) \widehat{\boldsymbol{f}}(O(-t) \boldsymbol{\xi}) \mathrm{d} t \\
& =\frac{1}{\omega} \int_{0}^{2 \pi} \Psi(\lambda, \boldsymbol{\xi}, t) O(t) \mathcal{F}[\boldsymbol{f}(O(-t) \cdot)](\boldsymbol{\xi}) \mathrm{d} t \tag{3.11}
\end{align*}
$$

where

$$
\Psi(\lambda, \boldsymbol{\xi}, t)=\frac{\mathrm{e}^{-t a(\xi) / \omega}}{\mathrm{e}^{-2 \pi a(\xi) / \omega}-1}
$$

In order to complete the proof, we shall need the next lemma.
Lemma 3.2. If $\lambda \notin \Lambda_{\gamma}^{\omega}$ then there exists a positive constant $c_{4}$ depending only on $\gamma$, $\omega$ and on the position of $\lambda$ in $\boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$ such that

$$
\begin{equation*}
\left|\mathrm{e}^{-2 \pi\left(\lambda+\mathrm{i} \gamma \xi_{1}+\nu s^{2}\right) / \omega}-1\right| \geq c_{4} \quad \text { for all } \xi \in \boldsymbol{R}^{3} \tag{3.12}
\end{equation*}
$$

Proof. The modulus of $\mathrm{e}^{-2 \pi a(\boldsymbol{\xi}) / \omega}-1$ is bounded below by a positive constant depending only on $\lambda, \gamma$ and $\omega$ (i.e., independent of $\boldsymbol{\xi}$ ) if the distance $d_{\lambda}(\boldsymbol{\xi}, k)$ of the complex number $-2 \pi a(\boldsymbol{\xi}) / \omega$ from $2 \pi \mathrm{i} k$ (for $k \in \boldsymbol{Z}$ ) is bounded below by another positive constant, independent of $\boldsymbol{\xi}$ and $k$. Let $\lambda=\alpha+\mathrm{i} \beta$ and $\rho^{2}=\xi_{2}^{2}+\xi_{3}^{2}$ so that

$$
\begin{aligned}
d_{\lambda}^{2}(\boldsymbol{\xi}, k) & =\frac{4 \pi^{2}}{\omega^{2}}\left|\alpha+\mathrm{i} \beta+\mathrm{i} \gamma \xi_{1}+\nu \xi_{1}^{2}+\nu \rho^{2}+\mathrm{i} k \omega\right|^{2} \\
& =\frac{4 \pi^{2}}{\omega^{2}}\left(\alpha+\nu \xi_{1}^{2}+\nu \rho^{2}\right)^{2}+\frac{4 \pi^{2}}{\omega^{2}}\left(\beta+\gamma \xi_{1}+k \omega\right)^{2} .
\end{aligned}
$$

Assume that $\min _{\xi \in \boldsymbol{R}^{3}} d_{\lambda}^{2}(\xi, k)=0$ for some fixed $k \in \boldsymbol{Z}$; the minimum exists because $\lim _{|\xi| \rightarrow \infty} d_{\lambda}^{2}(\xi, k)=\infty$. Then $\xi_{1}=-(\beta+k \omega) / \gamma$ and

$$
\alpha=-v \rho^{2}-\frac{v}{\gamma^{2}}(\beta+k \omega)^{2} \leq-\frac{v}{\gamma^{2}}(\beta+k \omega)^{2},
$$

but this inequality contradicts with $\lambda=\alpha+\mathrm{i} \beta \notin \Lambda_{\gamma}^{\omega}$. Hence

$$
D_{\lambda}(k):=\min _{\xi \in \boldsymbol{R}^{3}} d_{\lambda}(\boldsymbol{\xi}, k)>0
$$

for every $k \in \boldsymbol{Z}$. Moreover, the sequence $\left\{D_{\lambda}(k)\right\}_{k \in \boldsymbol{Z}}$ does not converge to zero as $k \rightarrow+\infty$ or $k \rightarrow-\infty$, because this would imply the existence of $\boldsymbol{\xi}(k) \in \boldsymbol{R}^{3}$ such that $d_{\lambda}(\boldsymbol{\xi}(k), k) \rightarrow 0$ as $k \rightarrow+\infty$ or $k \rightarrow-\infty$. Consequently $\beta+\gamma \xi_{1}(k)+k \omega \rightarrow 0,\left|\xi_{1}(k)\right| \rightarrow \infty$ and $\alpha+v \xi_{1}(k)^{2}+v \rho(k)^{2} \rightarrow \infty$, which is in contradiction with $d_{\lambda}(\boldsymbol{\xi}(k), k) \rightarrow 0$. Thus Lemma 3.2 is proved.

Continuation of the proof of Lemma 3.1. Lemma 3.2 implies that, for fixed $\lambda \in \boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$, the modulus of $\Psi(\lambda, \boldsymbol{\xi}, t)$ is bounded uniformly with respect to $t \in(0,2 \pi)$ and $\boldsymbol{\xi} \in \boldsymbol{R}^{3}$. Further, if $i \in\{2,3\}$ and $j \in\{1,2,3\}$, then

$$
\frac{\partial \Psi}{\partial \xi_{i}}=\frac{2 \nu \xi_{i}}{\omega} \frac{-t \mathrm{e}^{-t a(\xi) / \omega}\left(\mathrm{e}^{-2 \pi a(\xi) / \omega}-1\right)+2 \pi \mathrm{e}^{-(2 \pi+t) a(\xi) / \omega}}{\left[\mathrm{e}^{-2 \pi a(\xi) / \omega}-1\right]^{2}}
$$

and $\left|\xi_{i}\right|\left|\partial \Psi / \partial \xi_{i}\right|$ can be estimated as follows:

$$
\begin{aligned}
\left|\xi_{i}\right|\left|\frac{\partial \Psi}{\partial \xi_{i}}\right| & \leq\left|\frac{2 v s^{2}}{\omega} \frac{t \mathrm{e}^{-t a(\xi) / \omega}+(2 \pi-t) \mathrm{e}^{-(2 \pi+t) a(\xi) / \omega}}{\left[\mathrm{e}^{-2 \pi a(\xi) / \omega}-1\right]^{2}}\right| \\
& \leq C(\lambda)\left[s^{2} \frac{t}{\omega} \mathrm{e}^{-\nu s^{2} t / \omega}+s^{2} \frac{2 \pi+t}{\omega} \mathrm{e}^{-\nu s^{2}(2 \pi+t) / \omega}\right] .
\end{aligned}
$$

We observe that the right-hand side is less than or equal to a constant independent of $\boldsymbol{\xi}$ and $t$. We can similarly estimate $\left|\xi_{j}\right|\left|\partial \Psi / \partial \xi_{1}\right|$ and all other terms of the form

$$
\left|\xi_{1}\right|^{\kappa_{1}}\left|\xi_{2}\right|^{\kappa_{2}}\left|\xi_{3}\right|^{\kappa_{3}}\left|\frac{\partial^{\kappa_{1}+\kappa_{2}+\kappa_{3}} \Psi}{\partial \xi_{1}^{\kappa_{1}} \partial \xi_{2}^{\kappa_{2}} \partial \xi_{3}^{\kappa_{3}}}\right|, \quad \kappa_{1}, \kappa_{2}, \kappa_{3} \in\{0,1\} .
$$

Applying the inverse Fourier transform to (3.11), we arrive at the formula

$$
\boldsymbol{v}(\boldsymbol{x}):=\frac{1}{\omega} \int_{0}^{2 \pi} O(t) \mathcal{F}_{-1}[\Psi(\lambda, \boldsymbol{\xi}, t) \mathcal{F}[\boldsymbol{f}(O(-t) \cdot)](\boldsymbol{\xi})](\boldsymbol{x}) \mathrm{d} t
$$

Using Lizorkin's multiplier theorem (see e.g. [15, p. 375]) and the estimates of $\Psi$ and its derivatives discussed above, we derive the inequality

$$
\begin{equation*}
\left\|\mathcal{F}_{-1}[\Psi(\lambda, \boldsymbol{\xi}, t) \mathcal{F}[\boldsymbol{f}(O(-t) \cdot)](\boldsymbol{\xi})]\right\|_{q} \leq c_{5}\|\boldsymbol{f}\|_{q}, \quad t \in[0,2 \pi] \tag{3.13}
\end{equation*}
$$

where $c_{5}=c_{5}(\lambda, q)$. Then (3.11) and (3.13) imply that there exists $c_{3}>0$, independent of $f$ and $\boldsymbol{v}$, such that $\boldsymbol{v}$ satisfies the estimate (3.4).

Lemma 3.3. Suppose that $\lambda \in \boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$. Then $\lambda \in \rho\left(A_{\gamma}^{\omega}\right)$.
Proof. The estimates (2.3) and (3.4) in Lemmas 2.2 and 3.1 imply that the range of $A_{\gamma}^{\omega}-\lambda I$ is closed in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$ and that the operator is injective. By similar arguments, the same result does hold for its adjoint, $\left(A_{\gamma}^{\omega}\right)^{*}-\bar{\lambda} I$, on the dual space $L_{\sigma}^{q^{\prime}}\left(\boldsymbol{R}^{3}\right)$ of $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. Since $\left(A_{\gamma}^{\omega}\right)^{*}-\bar{\lambda} I$ is injective, we conclude that $R\left(A_{\gamma}^{\omega}-\lambda I\right)=L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. This proves that $\lambda \in \rho\left(A_{\gamma}^{\omega}\right)$.

Lemma 3.4. Let $1<q<\infty$. Then $\sigma_{\mathrm{p}}\left(A_{\gamma}^{\omega}\right)=\emptyset$.
Proof. Let $\lambda=\alpha+\mathrm{i} \beta \in \Lambda_{\gamma}^{\omega}, \alpha \leq 0, \beta \in \boldsymbol{R}$, and let $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ satisfy the equation $A_{\gamma}^{\omega} \boldsymbol{v}-\lambda \boldsymbol{v}=\mathbf{0}$. Applying the Fourier transform, we arrive at the identity

$$
\omega \partial_{\phi} \widehat{\boldsymbol{v}}-\left(\lambda+\mathrm{i} \gamma \xi_{1}+\nu|\boldsymbol{\xi}|^{2}\right) \widehat{\boldsymbol{v}}-\omega \times \widehat{\boldsymbol{v}}=\mathbf{0} .
$$

(We are using the same notation as in the proof of Lemma 3.1.) First we consider the simpler case when $1<q \leq 2$, in which $\widehat{\boldsymbol{v}}$ is a function from $L^{q^{\prime}}\left(\boldsymbol{R}^{3}\right)^{3}$. Denoting $\widehat{\boldsymbol{w}}\left(\rho, \phi, \xi_{1}\right)=$ $O(-\phi)^{T} \widehat{v}\left(\rho, \phi, \xi_{1}\right)$, we arrive, by analogy with (3.8), at the equation

$$
\begin{equation*}
\omega \partial_{\phi} \widehat{\boldsymbol{w}}-\left(\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}\right) \widehat{\boldsymbol{w}}=\mathbf{0}, \quad s=|\boldsymbol{\xi}| . \tag{3.14}
\end{equation*}
$$

Solving explicitly this ordinary differential equation, we obtain

$$
\widehat{\boldsymbol{w}}\left(\rho, \phi+2 \pi, \xi_{1}\right)=\widehat{\boldsymbol{w}}\left(\rho, \phi, \xi_{1}\right) \mathrm{e}^{2 \pi\left(\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}\right)}
$$

Since $\widehat{\boldsymbol{w}}$ is $2 \pi$-periodic in the variable $\phi$, we get the impossible condition $\alpha+\nu s^{2}=0$ for a.a. $\boldsymbol{\xi} \in \boldsymbol{R}^{3}$ unless $\widehat{\boldsymbol{w}}=\mathbf{0}$. Hence $\widehat{\boldsymbol{v}}=\mathbf{0}$ a.e. in $\boldsymbol{R}^{3}$ and also $\boldsymbol{v}=\mathbf{0}$ in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$.

In the general case $q>2$ let us again fix $\lambda \in \Lambda_{\gamma}^{\omega}, \lambda=\alpha+\mathrm{i} \beta+\mathrm{i} k \omega$ where $\alpha \leq-\nu \beta^{2} / \gamma^{2}$, $k \in \boldsymbol{Z}$. Consider $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and $\widehat{\boldsymbol{w}}\left(\rho, \phi, \xi_{1}\right)=O(\phi) \widehat{\boldsymbol{v}}\left(\rho, \phi, \xi_{1}\right)$ as above. Since the coefficients of $O(\phi)$ are either constant or $\cos \phi=\xi_{2} / \rho$ or $\sin \phi=\xi_{3} / \rho$, the function $\boldsymbol{w}$ (the inverse Fourier transform of $\widehat{\boldsymbol{w}}$ ) is defined by an application of 2D-Riesz transforms to $\boldsymbol{v}$. Hence $\boldsymbol{w} \in L^{q}\left(\boldsymbol{R}^{3}\right)^{3}$.

First we determine the support of tempered distributions $\widehat{\boldsymbol{w}}$ solving (3.14).
ASSERTION 3.1. Let $\widehat{\boldsymbol{w}} \in \mathcal{S}\left(\boldsymbol{R}^{3}\right)^{3}$ be a solution of (3.14). Then

$$
\operatorname{supp} \widehat{\boldsymbol{w}} \subset D:=\left\{\boldsymbol{\xi} \in \boldsymbol{R}^{3} ; \alpha+v s^{2}=0, \beta+\gamma \xi_{1} \in \omega \mathbf{Z}\right\}
$$

Proof of Assertion 3.1. Given $\boldsymbol{\psi} \in C_{0}^{\infty}\left(\boldsymbol{R}^{3} \backslash D\right)^{3}$, we solve the equation

$$
\begin{equation*}
-\omega \partial_{\phi} \boldsymbol{\Psi}-a(\xi) \boldsymbol{\Psi}=\boldsymbol{\psi} \tag{3.15}
\end{equation*}
$$

where $a(\boldsymbol{\xi})=\lambda+\mathrm{i} \gamma \xi_{1}+\nu s^{2}, s=|\boldsymbol{\xi}|$. Obviously, (3.15) yields the solution

$$
\boldsymbol{\Psi}(\phi)=\mathrm{e}^{-a(\xi) \phi / \omega}\left(\boldsymbol{\Psi}_{0}-\frac{1}{\omega} \int_{0}^{\phi} \mathrm{e}^{a(\xi) \phi^{\prime} / \omega} \boldsymbol{\psi}\left(\phi^{\prime}\right) \mathrm{d} \phi^{\prime}\right)
$$

when omitting the variables $\rho, \xi_{1}$ in $\boldsymbol{\Psi}$ and $\boldsymbol{\psi}$. Since $\boldsymbol{\Psi}$ is $2 \pi$-periodic in $\phi$, the initial value $\boldsymbol{\Psi}_{0}$ must satisfy the condition

$$
\left(\mathrm{e}^{-2 \pi a(\xi) / \omega}-1\right) \boldsymbol{\Psi}_{0}=\frac{1}{\omega} \mathrm{e}^{-2 \pi a(\boldsymbol{\xi}) / \omega} \int_{0}^{2 \pi} \mathrm{e}^{a(\xi) \phi^{\prime} / \omega} \boldsymbol{\psi}\left(\phi^{\prime}\right) \mathrm{d} \phi^{\prime} .
$$

This equation is uniquely solvable for the unknown $\boldsymbol{\Psi}_{0}$ if

$$
2 \pi a(\boldsymbol{\xi}) / \omega \notin 2 \pi \mathrm{i} \boldsymbol{Z} \Leftrightarrow \lambda+\mathrm{i} \gamma \xi_{1}+v s^{2} \neq k \mathrm{i} \omega \quad \text { for all } k \in \boldsymbol{Z} \Leftrightarrow \boldsymbol{\xi} \notin D .
$$

For $\boldsymbol{\xi} \notin D$, we get the unique solution of (3.15)

$$
\boldsymbol{\Psi}(\phi)=\frac{1}{\omega} \frac{1}{1-\mathrm{e}^{2 \pi a(\xi) / \omega}} \int_{0}^{2 \pi} \mathrm{e}^{a(\xi) \phi^{\prime} / \omega} \boldsymbol{\psi}\left(\phi^{\prime}+\phi\right) \mathrm{d} \phi^{\prime}
$$

in view of (3.10) with $\omega$ replaced by $-\omega$. Since supp $\boldsymbol{\psi} \subset \boldsymbol{R}^{3} \backslash D$, we obtain that $\boldsymbol{\Psi} \in$ $C_{0}^{\infty}\left(\boldsymbol{R}^{3} \backslash D\right)^{3}$.

Now we use (3.14) to get that for all $\boldsymbol{\psi} \in C_{0}^{\infty}\left(\boldsymbol{R}^{3} \backslash D\right)^{3}$

$$
\langle\widehat{\boldsymbol{w}}, \boldsymbol{\psi}\rangle=\left\langle\widehat{\boldsymbol{w}},-\omega \partial_{\phi} \boldsymbol{\Psi}-a(\boldsymbol{\xi}) \boldsymbol{\Psi}\right\rangle=\left\langle\omega \partial_{\phi} \widehat{\boldsymbol{w}}-\left(\lambda+\mathrm{i} \gamma \xi_{1}+v s^{2}\right) \widehat{\boldsymbol{w}}, \boldsymbol{\Psi}\right\rangle=0 .
$$

This identity proves that supp $\widehat{\boldsymbol{w}} \subset D$.
Continuation of the proof of Lemma 3.4. The set $D$ can be written as the union of finitely many disjoint sets of type

$$
D_{k}:=\left\{\boldsymbol{\xi} \in \boldsymbol{R}^{3} ; \alpha+v s^{2}=0, \beta+\gamma \xi_{1}=\omega k\right\}, \quad k \in \boldsymbol{Z}
$$

each non-void set $D_{k}$ defines a circle in $\boldsymbol{R}^{3}$ parallel to the $\xi_{2} \xi_{3}$-plane with center $\left(\xi_{1}^{(k)}, 0,0\right)$ where $\xi_{1}^{(k)}=(\omega k-\beta) / \gamma$. At least the set $D_{0}$ is non-void since $\lambda \in \Lambda_{\gamma}^{\omega}$. Using a suitable partition of unity with respect to the variable $\xi_{1}$, we may write $\widehat{\boldsymbol{w}} \in \mathcal{S}^{\prime}\left(\boldsymbol{R}^{3}\right)^{3}$, a solution of
(3.14), as a finite linear combination of tempered distributions $\widehat{\boldsymbol{w}}_{k}$ with supp $\widehat{\boldsymbol{w}}_{k} \subset D_{k}$. Since $\boldsymbol{w}$ is in $L^{q}\left(\boldsymbol{R}^{3}\right)^{3}, \boldsymbol{w}_{k}$ is also in $L^{q}\left(\boldsymbol{R}^{3}\right)^{3}$ for each $k$. Moreover, since $\boldsymbol{w}$ is in $C^{\infty}\left(\boldsymbol{R}^{3}\right)^{3}$ by elliptic regularity theory, $\boldsymbol{w}_{k}$ is also in $C^{\infty}\left(\boldsymbol{R}^{3}\right)^{3}$. We shall further need the assertion:

AsSERTION 3.2. Let $h \in \mathcal{S}^{\prime}\left(\boldsymbol{R}^{3}\right) \cap C^{\infty}\left(\boldsymbol{R}^{3}\right)$ satisfy supp $\widehat{h} \subset\{0\} \times \boldsymbol{R}^{2}$. Then $h$ is a polynomial with respect to the variable $x_{1}$.

Proof of Assertion 3.2. Given $h \in \mathcal{S}^{\prime}\left(\boldsymbol{R}^{3}\right)$, we find $N \in \boldsymbol{N}$ and a constant $c>0$ such that

$$
|\langle\widehat{h}, \psi\rangle| \leq c\left\|(1+|\cdot|)^{N} \sum_{0 \leq|\mu| \leq N}\left|\partial^{\mu} \psi\right|\right\|_{\infty}
$$

for all test functions $\psi \in \mathcal{S}\left(\boldsymbol{R}^{3}\right)$; here $\mu \in \boldsymbol{N}^{3}$ denotes multi-indices and $\|\cdot\|_{\infty}$ the supremum norm for functions on $\boldsymbol{R}^{3}$. Let $\eta \in C_{0}^{\infty}(\boldsymbol{R})$ be a cut-off function such that $\eta\left(\xi_{1}\right)=1$ for $\xi_{1}$ in a neighborhood of 0 , and let $\eta_{\varepsilon}\left(\xi_{1}\right)=\eta\left(\xi_{1} / \varepsilon\right)$ for $\varepsilon>0$. Note that $\left\|\partial_{1}^{m} \eta_{\varepsilon}\right\|_{\infty} \leq c_{m} \varepsilon^{-m}$ for all $m \in \boldsymbol{N}$. Since supp $\widehat{h} \subset\{0\} \times \boldsymbol{R}^{2}$, we get that $\langle\widehat{h}, \psi\rangle=\left\langle\widehat{h}, \eta_{\varepsilon} \psi\right\rangle$ for all $\varepsilon>0$. Hence, with a constant $c>0$ independent of $\varepsilon \in(0,1)$,

$$
\left|\left\langle\xi_{1}^{N+1} \widehat{h}, \psi\right\rangle\right|=\left|\left\langle\widehat{h}, \xi_{1}^{N+1} \eta_{\varepsilon} \psi\right\rangle\right| \leq c \varepsilon\left\|(1+|\cdot|)^{N} \sum_{0 \leq|\mu| \leq N}\left|\partial^{\mu} \psi\right|\right\|_{\infty}
$$

Consequently, $\xi_{1}^{N+1} \widehat{h}=0$. Now we conclude that $\partial_{1}^{N+1} h=0$ and that $h \in C^{\infty}\left(\boldsymbol{R}^{3}\right)$ is a polynomial with respect to the variable $x_{1}$ of order at most $N$.

Completion of the proof of Lemma 3.4. To complete the proof of Lemma 3.4, let $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ satisfy the equation $A_{\gamma}^{\omega} \boldsymbol{v}-\lambda \boldsymbol{v}=\mathbf{0}$. By the above arguments, it suffices to show that each $\boldsymbol{w}_{k}$ in the partition of $\boldsymbol{w}$ vanishes. Here $\boldsymbol{w}_{k} \in L^{q}\left(\boldsymbol{R}^{3}\right)^{3} \cap C^{\infty}\left(\boldsymbol{R}^{3}\right)^{3}$ has the property supp $\widehat{\boldsymbol{w}}_{k} \subset D_{k}$, i.e., supp $\widehat{\boldsymbol{w}}_{k} \subset\left\{\xi_{1}^{(k)}\right\} \times \boldsymbol{R}^{2}$ with $\xi_{1}^{(k)}=(\omega k-\beta) / \gamma$. Then the function $x \mapsto \mathrm{e}^{\mathrm{i} x_{1} \xi_{1}^{(k)}} \boldsymbol{w}_{k}(\boldsymbol{x})$, which satisfies the assumptions of Assertion 3.2, is a polynomial with respect to the variable $x_{1}$. Since this function is contained in $L^{q}\left(\boldsymbol{R}^{3}\right)^{3}$, it must vanish identically. Now we proved that $\boldsymbol{w}=\mathbf{0}$ and also $\boldsymbol{v}=\mathbf{0}$.

Lemma 3.5. Let $1<q<\infty$. Then $\sigma_{\mathrm{r}}\left(A_{\gamma}^{\omega}\right)=\emptyset$.
Proof. Lemma 3.4 and duality arguments yield the assertion.
Lemma 3.6. Let $1<q<\infty$. Then

$$
\sigma\left(A_{\gamma}^{\omega}\right)=\sigma_{\mathrm{c}}\left(A_{\gamma}^{\omega}\right)=\sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)=\Lambda_{\gamma}^{\omega} .
$$

PROOF. Lemmas 3.3 through 3.5 imply that $\sigma\left(A_{\gamma}^{\omega}\right)=\sigma_{\mathrm{c}}\left(A_{\gamma}^{\omega}\right) \subset \Lambda_{\gamma}^{\omega}$. Thus, we need to show the opposite inclusion, i.e., $\Lambda_{\gamma}^{\omega} \subset \sigma_{\mathrm{c}}\left(A_{\gamma}^{\omega}\right)$.

Suppose at first that $\lambda \in\left(\Lambda_{\gamma}^{\omega}\right)^{\circ}$, the interior of $\Lambda_{\gamma}^{\omega}$. It means that there exist $\alpha, \beta \in \boldsymbol{R}$ and $k \in Z$ such that $\lambda=\alpha+\mathrm{i} \beta+\mathrm{i} k \omega$ and $\alpha<-\nu \beta^{2} / \gamma^{2}$. The number $\alpha$ can be written in the form $\alpha=\alpha_{1}+\alpha_{2}$ where $\alpha_{1}=-\nu \beta^{2} / \gamma^{2}$ and $\alpha_{2}<0$. Assume that $k \neq 0$. The procedure in the case $k=0$ would be analogous.

We shall explicitly define a sequence of functions $\left\{\boldsymbol{v}^{n}\right\}$ in $D\left(A_{\gamma}^{\omega}\right)$ such that $\left\|\boldsymbol{v}^{n}\right\|_{q}=1$ and $\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) v^{n}\right\|_{q} \rightarrow 0$ as $n \rightarrow \infty$. It will imply that nul' $\left(A_{\gamma}^{\omega}-\lambda I\right)>0$.

Let us denote by $v_{1}^{n}, v_{r}^{n}$ and $v_{\varphi}^{n}$ the cylindrical components of $\boldsymbol{v}^{n}$. Put

$$
\begin{aligned}
& v_{1}^{n}\left(x_{1}, r, \varphi\right):=0 \\
& v_{r}^{n}\left(x_{1}, r, \varphi\right):=\kappa_{n} U^{n}\left(x_{1}\right) V^{n}(r) \mathrm{e}^{\mathrm{i} k \varphi}, \\
& v_{\varphi}^{n}\left(x_{1}, r, \varphi\right):=-\frac{1}{\mathrm{i} k} \partial_{r}\left[r v_{r}^{n}\left(x_{1}, r, \varphi\right)\right]=-\frac{1}{\mathrm{i} k} \kappa_{n} U^{n}\left(x_{1}\right)\left[V^{n}(r)+r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right] \mathrm{e}^{\mathrm{i} k \varphi} .
\end{aligned}
$$

Here the function $U^{n}$ has the form

$$
U^{n}\left(x_{1}\right):=\eta_{1}^{n}\left(x_{1}\right) Y\left(x_{1}\right), \quad \text { where } \quad Y\left(x_{1}\right)=\mathrm{e}^{\mathrm{i} a x_{1}}, a=-\frac{\beta}{\gamma}
$$

and where $\eta_{1}^{n}$ is an infinitely differentiable function on $(-\infty,+\infty)$ such that $0 \leq \eta_{1}^{n} \leq 1$,

$$
\eta_{1}^{n}\left(x_{1}\right)= \begin{cases}0 & \text { for } x_{1} \leq-n-n^{2} \text { and } n+n^{2} \leq x_{1}, \\ 1 & \text { for }-n^{2} \leq x_{1} \leq n^{2}\end{cases}
$$

The identity $\alpha_{1}=-\nu \beta^{2} / \gamma^{2}$ guarantees that the characteristic equation $\nu \zeta^{2}-\gamma \zeta-\left(\alpha_{1}+\mathrm{i} \beta\right)=$ 0 , corresponding to the equation (3.16) below, has the root $\zeta_{1}=\mathrm{i} a$. Thus, the function $Y$ is a bounded non-trivial solution of the ordinary differential equation

$$
\begin{equation*}
v Y^{\prime \prime}\left(x_{1}\right)-\gamma Y^{\prime}\left(x_{1}\right)-\left(\alpha_{1}+\mathrm{i} \beta\right) Y\left(x_{1}\right)=0 \tag{3.16}
\end{equation*}
$$

in the interval $(-\infty,+\infty)$. The function $V^{n}$ is supposed to have the form

$$
V^{n}(r):=\eta_{2}^{n}(r) \mathrm{e}^{\mathrm{i} b r}, \quad b=\sqrt{-\frac{\alpha_{2}}{v}},
$$

where $\eta_{2}^{n}$ is an infinitely differentiable function on $[0,+\infty)$ such that $0 \leq \eta_{2}^{n} \leq 1$ and

$$
\eta_{2}^{n}(r)= \begin{cases}0 & \text { for } 0 \leq r \leq n \text { and } 3 n+n^{2} \leq r, \\ 1 & \text { for } 2 n \leq r \leq 2 n+n^{2}\end{cases}
$$

Both the functions $\eta_{1}^{n}$ and $\eta_{2}^{n}$ can be chosen so that their derivatives are of the order $1 / n$. The definition of $V^{n}$ guarantees that it satisfies

$$
\begin{equation*}
v \frac{\mathrm{~d}^{2}}{\mathrm{~d} r^{2}} V^{n}(r)-\alpha_{2} V^{n}(r)=0 \tag{3.17}
\end{equation*}
$$

for $2 n<r<2 n+n^{2}$. Finally, the constant $\kappa_{n}$ is chosen so that $\left\|\boldsymbol{v}^{n}\right\|_{q}=1$.
One can easily check that $\boldsymbol{v}^{n}$ satisfies the condition of incompressibility

$$
\nabla \cdot \boldsymbol{v}^{n} \equiv \partial_{1} v_{1}^{n}+\frac{1}{r} \partial_{r}\left(r v_{r}^{n}\right)+\frac{1}{r} \partial_{\varphi} v_{\varphi}^{n}=0
$$

The support of $\boldsymbol{v}^{n}$ is a subset of

$$
\begin{equation*}
S^{n}:=\left\{\boldsymbol{x}=\left[x_{1}, r, \varphi\right] \in \boldsymbol{R}^{3} ;\left|x_{1}\right| \leq n+n^{2}, n \leq r \leq 3 n+n^{2}, 0 \leq \varphi<2 \pi\right\} . \tag{3.18}
\end{equation*}
$$

Considering the norm of $\boldsymbol{v}^{n}$, we observe that for large $n$ the decisive contribution comes from the integral of $\left|v_{\varphi}^{n}\right|^{q}$, namely of its part $\left|(-1 / \mathrm{i} k) \kappa_{n} U^{n} r\left(\mathrm{~d} V^{n} / \mathrm{d} r\right) \mathrm{e}^{\mathrm{i} k \varphi}\right|^{q}$, on the region

$$
S_{0}^{n}:=\left\{\boldsymbol{x}=\left[x_{1}, r, \varphi\right] \in \boldsymbol{R}^{3} ;\left|x_{1}\right|<n^{2}, 2 n<r<2 n+n^{2}, 0<\varphi<2 \pi\right\} .
$$

The integrals of all other parts on other regions are of a lower order in $n$. Calculating the integral of $\left|(-1 / \mathrm{i} k) \kappa_{n} U^{n} r\left(\mathrm{~d} V^{n} / \mathrm{d} r\right) \mathrm{e}^{\mathrm{i} k \varphi}\right|^{q}$ on the domain $S_{0}^{n}$, we obtain

$$
\begin{aligned}
\int_{-n^{2}}^{n^{2}} & \int_{2 n}^{2 n+n^{2}} \int_{0}^{2 \pi}\left|\frac{\kappa_{n}}{\mathrm{i} k} U^{n}\left(x_{1}\right) r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right|^{q} 2 r \mathrm{~d} \varphi \mathrm{~d} r \mathrm{~d} x_{1} \\
& =2 \pi \frac{\kappa_{n}^{q}}{|k|^{q}} \int_{-n^{2}}^{n^{2}}\left|U^{n}\left(x_{1}\right)\right|^{q} \mathrm{~d} x_{1} \int_{2 n}^{2 n+n^{2}} r^{q+1}\left|\frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right|^{q} \mathrm{~d} r \\
& =2 \pi \frac{\kappa_{n}^{q}}{|k|^{q}} 2 n^{2} \frac{b^{q}}{q+2}\left(\left(2 n+n^{2}\right)^{q+2}-(2 n)^{q+2}\right)
\end{aligned}
$$

Here we have used the equalities $\eta_{1}^{n}\left(x_{1}\right)=\eta_{2}^{n}(r)=1$, hence $\left|U^{n}\left(x_{1}\right)\right|=\left|V^{n}(r)\right|=1$ for $\left(x_{1}, r, \varphi\right) \in S_{0}^{n}$. Thus, there exist $n_{0} \in \boldsymbol{N}$ and positive constants $c_{6}$ and $c_{7}$ (independent of $n$ ) such that

$$
\begin{equation*}
\frac{c_{6}}{n^{2+6 / q}} \leq \kappa_{n} \leq \frac{c_{7}}{n^{2+6 / q}} \quad \text { for all } n \in \boldsymbol{N}, \text { with } n \geq n_{0} . \tag{3.19}
\end{equation*}
$$

We can naturally use the form (3.1) of the operator $A_{\gamma}^{\omega}$. Moreover, if we identify $\boldsymbol{v}$ with the triplet of its cylindrical coordinates, i.e., $\boldsymbol{v} \triangleq\left[v_{1}, v_{r}, v_{\varphi}\right]^{T}$, then we can verify that

$$
A_{\gamma}^{\omega} \boldsymbol{v}=v \Delta \boldsymbol{v}+\omega \partial_{\varphi} \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}
$$

(see [14]). Hence

$$
\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{v}^{n}=\left(v \Delta+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right)\left[\begin{array}{c}
0 \\
-\frac{1}{\mathrm{i} k} \kappa_{n} U^{n}\left(x_{1}\right)\left[V^{n}(r)+r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right]
\end{array} \mathrm{e}^{\mathrm{i} k \varphi} .\right] .
$$

Calculating now the norm of $\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{v}^{n}$ in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$, we observe that the contributions coming from $\boldsymbol{R}^{3} \backslash S_{0}^{n}$ tend to zero as $n \rightarrow \infty$ because they represent $q$-roots of integrals of functions bounded by $C \kappa_{n}^{q} r^{q}$ on $S^{n} \backslash S_{0}^{n}$. Due to (3.19), this contribution is of the order $n^{-1 / q}$. Concerning the integral on $S_{0}^{n}$, the decisive part again comes from $\left(\nu \Delta+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right) v_{\varphi}^{n}$, namely from $\left(\nu \Delta+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right)$ applied to the term $(-1 / \mathrm{i} k) \kappa_{n} U^{n} r\left(\mathrm{~d} V^{n} / \mathrm{d} r\right) \mathrm{e}^{\mathrm{i} k \varphi}$ because of the factor $r$ inside this term. By the equality $\lambda=\alpha_{1}+\alpha_{2}+\mathrm{i} \beta+\mathrm{i} k \omega$ and (3.16) and (3.17), we have

$$
\begin{aligned}
& \left(v \Delta+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right)\left(\frac{\kappa_{n}}{\mathrm{i} k} U^{n}\left(x_{1}\right) r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r} \mathrm{e}^{\mathrm{i} k \varphi}\right) \\
& \quad=\frac{\kappa_{n}}{\mathrm{i} k}\left(v \partial_{1}^{2}+v \partial_{r}^{2}+\frac{v}{r} \partial_{r}+\frac{v}{r^{2}} \partial_{\varphi}^{2}+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right)\left(U^{n}\left(x_{1}\right) r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r} \mathrm{e}^{\mathrm{i} k \varphi}\right)
\end{aligned}
$$

$$
\begin{aligned}
= & \frac{\kappa_{n}}{\mathrm{i} k} U^{n}\left(x_{1}\right)\left(v \frac{\mathrm{~d}^{2}}{\mathrm{~d} r^{2}}+\frac{\nu}{r} \frac{\mathrm{~d}}{\mathrm{~d} r}-\alpha_{2} I\right)\left[r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right] \mathrm{e}^{\mathrm{i} k \varphi} \\
& +\frac{\kappa_{n}}{\mathrm{i} k}\left(\nu Y^{\prime \prime}\left(x_{1}\right)-\gamma Y^{\prime}\left(x_{1}\right)-\left[\alpha_{1}+\mathrm{i} \beta\right] Y\left(x_{1}\right)\right)\left[r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right] \mathrm{e}^{\mathrm{i} k \varphi} \\
& -\frac{\kappa_{n}}{\mathrm{i} k} \frac{\nu k^{2}}{r^{2}} U^{n}\left(x_{1}\right) r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r} \mathrm{e}^{\mathrm{i} k \varphi} \\
= & \frac{\kappa_{n}}{\mathrm{i} k}\left\{U^{n}\left(x_{1}\right) r \frac{\mathrm{~d}}{\mathrm{~d} r}\left[v \frac{\mathrm{~d}^{2} V^{n}(r)}{\mathrm{d} r^{2}}-\alpha_{2} V^{n}(r)\right]+U^{n}\left(x_{1}\right) 2 v \frac{\mathrm{~d}^{2} V^{n}(r)}{\mathrm{d} r^{2}}\right. \\
& \left.+U^{n}\left(x_{1}\right) \frac{\nu}{r} \frac{\mathrm{~d}}{\mathrm{~d} r}\left[r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right]-\frac{\kappa_{n}}{\mathrm{i} k} U^{n}\left(x_{1}\right) \frac{\nu k^{2}}{r} \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r}\right\} \mathrm{e}^{\mathrm{i} k \varphi} \\
= & \frac{\nu \kappa_{n}}{\mathrm{i} k}\left(-3 b^{2}+\frac{\mathrm{i} b}{r}-\frac{k^{2} \mathrm{i} b}{r}\right) \mathrm{e}^{\mathrm{i}\left(a x_{1}+b r\right)} \mathrm{e}^{\mathrm{i} k \varphi},
\end{aligned}
$$

where, in the last step, we used the simple forms of the functions $U^{n}$ and $V^{n}$ on $S_{0}^{n}$, i.e., $U^{n}\left(x_{1}\right)=\mathrm{e}^{\mathrm{i} a x_{1}}$ and $V^{n}(r)=\mathrm{e}^{\mathrm{i} b r}$. Hence

$$
\begin{aligned}
& {\left[\int_{-n^{2}}^{n^{2}} \int_{2 n}^{2 n+n^{2}} \int_{0}^{2 \pi}\left|\left(\nu \Delta+\omega \partial_{\varphi}-\gamma \partial_{1}-\lambda I\right)\left(\frac{\kappa_{n}}{\mathrm{i} k} U^{n}\left(x_{1}\right) r \frac{\mathrm{~d} V^{n}(r)}{\mathrm{d} r} \mathrm{e}^{\mathrm{i} k \varphi}\right)\right|^{q} r \mathrm{~d} \varphi \mathrm{~d} r \mathrm{~d} x_{1}\right]^{1 / q}} \\
& \quad \leq C(\nu, k, b) \kappa_{n}\left[\int_{-n^{2}}^{n^{2}} \int_{2 n}^{2 n+n^{2}} r \mathrm{~d} r \mathrm{~d} x_{1}\right]^{1 / q} \\
& \quad=C(v, k, b) \kappa_{n}\left\{2 n^{2}\left[\left(2 n+n^{2}\right)^{2}-(2 n)^{2}\right]\right\}^{1 / q}
\end{aligned}
$$

The last term tends to zero as $n \rightarrow \infty$ by (3.19). In this way, we are led to the convergence $\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{v}^{n}\right\|_{q} \rightarrow 0$ as $n \rightarrow \infty$.

We have proved that $\operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)>0$. Hence nul $\left(A_{\gamma}^{\omega}-\lambda I\right) \neq \operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)$. It means that the range $R\left(A_{\gamma}^{\omega}-\lambda I\right)$ is not closed in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. Consequently, $\operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)=$ $\operatorname{def}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)=+\infty$ and $\lambda \in \sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$. Since $\sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$ is closed, $\Lambda_{\gamma}^{\omega}$ is a subset of $\sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$. By Lemmas 3.4 and 3.5, we have the inclusion $\Lambda_{\gamma}^{\omega} \subset \sigma_{\mathrm{c}}\left(A_{\gamma}^{\omega}\right)$.

If $q=2$, then $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right) \equiv L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$ is a Hilbert space and it is natural to ask whether the operator $A_{\gamma}^{\omega}$ is normal. The answer is given by the next lemma.

Lemma 3.7. Let $q=2$. Then $A_{\gamma}^{\omega}$ is a normal operator in $L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$.
Proof. Using the cylindrical coordinates $\left(x_{1}, r, \varphi\right)$ as in the proof of Lemma 3.1, we can express the operators $A_{\gamma}^{\omega}$ and $\left(A_{\gamma}^{\omega}\right)^{*}$ in accordance with (3.1), (3.2) and (3.6) as

$$
\begin{aligned}
A_{\gamma}^{\omega} \boldsymbol{v} & =v \Delta \boldsymbol{v}+\omega \partial_{\varphi} \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}, \\
\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{v} & =v \Delta \boldsymbol{v}-\omega \partial_{\varphi} \boldsymbol{v}+\boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v},
\end{aligned}
$$

where $D\left(A_{\gamma}^{\omega}\right)=D\left(\left(A_{\gamma}^{\omega}\right)^{*}\right)$. We need to show that $A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*}=\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}$, i.e.,
(a) $D\left(A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*}\right)=D\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}\right)$,
(b) $A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}=\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega} \boldsymbol{u}$ for all $\boldsymbol{u} \in D\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}\right)$.

Let us begin with part (a). Suppose that $\boldsymbol{u} \in D\left(A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*}\right)$, i.e., $\boldsymbol{u} \in D\left(\left(A_{\gamma}^{\omega}\right)^{*}\right)$ and $\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u} \in D\left(A_{\gamma}^{\omega}\right)$. In order to show that $\boldsymbol{u} \in D\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}\right)$, we treat the scalar product $\left(A_{\gamma}^{\omega} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}$ of $A_{\gamma}^{\omega} \boldsymbol{u}$ and $A_{\gamma}^{\omega} \boldsymbol{v}$ in $L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$ for $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ as follows:

$$
\begin{align*}
\left(A_{\gamma}^{\omega} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}= & \left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}+2\left(\omega \partial_{\varphi} \boldsymbol{u}-\boldsymbol{\omega} \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2} \\
= & \left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}+2\left(\omega \partial_{\varphi} \boldsymbol{u}-\omega \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}, v \Delta \boldsymbol{v}\right)_{2}  \tag{3.20}\\
& +2\left(\omega \partial_{\varphi} \boldsymbol{u}-\omega \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}, \omega \partial_{\varphi} \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}\right)_{2} .
\end{align*}
$$

Let us first assume that $\boldsymbol{v}$ has a compact support. Then $\left(\partial_{\varphi} \boldsymbol{u}, \Delta \boldsymbol{v}\right)_{2}=-\left(\Delta \boldsymbol{u}, \partial_{\varphi} \boldsymbol{v}\right)_{2}$ and

$$
\left(-\boldsymbol{\omega} \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}, v \Delta \boldsymbol{v}\right)_{2}=\left(v \Delta \boldsymbol{u}, \boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v}\right)_{2} .
$$

Substituting these identities into (3.20), we obtain

$$
\begin{align*}
\left(A_{\gamma}^{\omega} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}= & \left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}-2\left(v \Delta \boldsymbol{u}, \omega \partial_{\varphi} \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}\right)_{2} \\
& +2\left(\omega \partial_{\varphi} \boldsymbol{u}-\omega \times \boldsymbol{u}-\gamma \partial_{1} \boldsymbol{u}, \omega \partial_{\varphi} \boldsymbol{v}-\boldsymbol{\omega} \times \boldsymbol{v}-\gamma \partial_{1} \boldsymbol{v}\right)_{2} \\
= & \left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}+2\left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u},-\omega \partial_{\varphi} \boldsymbol{v}+\boldsymbol{\omega} \times \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{v}\right)_{2}  \tag{3.21}\\
= & \left(\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u},\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{v}\right)_{2}=\left(A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, \boldsymbol{v}\right)_{2} .
\end{align*}
$$

In fact, (3.21) holds for all $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ because by Lemma 2.3 the set $D_{0}\left(A_{\gamma}^{\omega}\right)=\{\boldsymbol{v} \in$ $D\left(A_{\gamma}^{\omega}\right) ; \boldsymbol{v}$ has a compact support in $\left.\boldsymbol{R}^{3}\right\}$ is a core of $A_{\gamma}^{\omega}$. Now, (3.21) shows that for fixed $\boldsymbol{u},\left(A_{\gamma}^{\omega} \boldsymbol{u}, A_{\gamma}^{\omega} \boldsymbol{v}\right)_{2}$ can be extended to a continuous linear functional of $\boldsymbol{v} \in L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$. Thus, $\boldsymbol{u} \in D\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}\right)$.

We have proved the inclusion $D\left(A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*}\right) \subset D\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega}\right)$. The opposite inclusion can be proved in the same way.

Concerning part (b), (3.21) implies that

$$
\left(\left(A_{\gamma}^{\omega}\right)^{*} A_{\gamma}^{\omega} \boldsymbol{u}, \boldsymbol{v}\right)_{2}=\left(A_{\gamma}^{\omega}\left(A_{\gamma}^{\omega}\right)^{*} \boldsymbol{u}, \boldsymbol{v}\right)_{2}
$$

for all $\boldsymbol{v} \in D\left(A_{\gamma}^{\omega}\right)$ and even for all $\boldsymbol{v} \in L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$ by the density of $D\left(A_{\gamma}^{\omega}\right)$ in $L_{\sigma}^{2}\left(\boldsymbol{R}^{3}\right)$. Hence the operators $A_{\gamma}^{\omega}$ and $\left(A_{\gamma}^{\omega}\right)^{*}$ commute.

Theorem 1.2 resumes the results of Lemmas 3.3 through 3.7.
4. The case of an exterior domain $\Omega$. In this section, we assume that $\Omega \subset \boldsymbol{R}^{3}$ is an exterior domain, different from $\boldsymbol{R}^{3}$, with boundary of class $C^{1,1}$.

Lemma 4.1. The inclusion $\Lambda_{\gamma}^{\omega} \subset \sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)$ holds.
Proof. We prove in almost the same way as in the proof of Lemma 3.6 that if $\lambda \in$ $\left(\Lambda_{\gamma}^{\omega}\right)^{\circ}$ then $\operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)>0$. Unfortunately, since an analogue to Lemma 3.4 is not available, we cannot deduce directly from this inequality that $\lambda \in \sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$ as in the proof of Lemma 3.6. However, we will use that supp $\boldsymbol{v}^{n} \subset S^{n}$ (see (3.18)). Thus, there exists $n_{0} \in \boldsymbol{N}$ so large that $\boldsymbol{v}^{n}$ belongs to the domain of $A_{\gamma}^{\omega}$ (as an operator in $L_{\sigma}^{q}(\Omega)$ ) for $n \geq n_{0}$. We observe that any subsequence $\left\{S^{k_{n}}\right\}$ of $\left\{S^{n}\right\}$ has the intersection property $\bigcap_{n=1}^{\infty} S^{k_{n}}=\emptyset$. It implies that the sequence $\left\{\boldsymbol{v}^{n}\right\}_{n \geq n_{0}}$ is not compact in $L_{\sigma}^{q}(\Omega)$. Consequently, $\operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)=$
$\infty$. We can prove in the same way that $\operatorname{nul}^{\prime}\left(\left(A_{\gamma}^{\omega}\right)^{*}-\bar{\lambda} I\right)=\infty$. Hence $\operatorname{def}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)=\infty$, and the operator $A_{\gamma}^{\omega}-\lambda I$ is not semi-Fredholm. Thus, $\lambda \in \sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right)$. The inclusion $\Lambda_{\gamma}^{\omega} \subset$ $\sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$ now follows from the closedness of $\sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$.

Lemma 4.2. The inclusion $\sigma_{\mathrm{ess}}\left(A_{\gamma}^{\omega}\right) \subset \Lambda_{\gamma}^{\omega}$ holds.
Proof. Let $\lambda \in \sigma_{\text {ess }}\left(A_{\gamma}^{\omega}\right)$. Then $\operatorname{nul}^{\prime}\left(A_{\gamma}^{\omega}-\lambda I\right)=\infty$. This information enables us to construct, by mathematical induction, a sequence $\left\{\boldsymbol{u}^{n}\right\}$ in $D\left(A_{\gamma}^{\omega}\right)$ satisfying $\left\|\boldsymbol{u}^{n}\right\|_{q}=1$, $\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{u}^{n}\right\|_{q} \rightarrow 0$ as $n \rightarrow \infty$ and

$$
\begin{equation*}
\operatorname{dist}\left(\boldsymbol{u}^{n} ; \mathcal{L}_{n-1}\right)=1 \tag{4.1}
\end{equation*}
$$

for all $n \in N$, where $\mathcal{L}_{n-1}$ denotes the linear hull of the functions $\boldsymbol{u}^{1}, \ldots, \boldsymbol{u}^{n-1}$ : Suppose that we have already constructed $\boldsymbol{u}^{1}, \ldots, \boldsymbol{u}^{k}$ satisfying $\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{u}^{j}\right\|_{q} \leq 1 / j$ for $j=1, \ldots, k$ and (4.1) for all $n=1, \ldots, k$. To $\varepsilon_{k+1}=1 /(k+1)$ there exists an infinite dimensional linear manifold $M_{k+1}$ in $D\left(A_{\gamma}^{\omega}\right)$ such that $\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{u}\right\|_{q} \leq \varepsilon_{k+1}$ for all $\boldsymbol{u} \in M_{k+1}$. Then due to [25, Lemma IV.2.3], we find $\boldsymbol{u}^{k+1} \in M_{k+1}$ such that $\left\|\boldsymbol{u}^{k+1}\right\|_{q}=1$ and $\operatorname{dist}\left(\boldsymbol{u}^{k+1} ; \mathcal{L}_{k}\right)=1$. The sequence $\left\{\boldsymbol{u}^{n}\right\}$ satisfies

$$
\begin{equation*}
\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{u}^{n}\right\|_{q} \leq \frac{1}{n} \quad \text { for all } n \in N \tag{4.2}
\end{equation*}
$$

Denote $\boldsymbol{f}^{n}:=\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{u}^{n}$. Lemma 2.2 yields the estimates

$$
\begin{align*}
\left\|\boldsymbol{u}^{n}\right\|_{2, q}+\left\|(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{u}^{n}\right\|_{q} & \leq c_{1}\left\|\boldsymbol{f}^{n}\right\|_{q}+\left(c_{2}+c_{1}|\lambda|\right)\left\|\boldsymbol{u}^{n}\right\|_{q} \\
& \leq c_{1}+\left(c_{2}+c_{1}|\lambda|\right):=c_{8} \tag{4.3}
\end{align*}
$$

with a constant $c_{8}>0$ independent of $n \in N$. Furthermore, there exists $\nabla p^{n} \in L^{q}(\Omega)^{3}$ such that

$$
\begin{equation*}
\nu \Delta \boldsymbol{u}^{n}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{u}^{n}-\boldsymbol{\omega} \times \boldsymbol{u}^{n}-\gamma \partial_{1} \boldsymbol{v}^{n}-\lambda \boldsymbol{u}^{n}-\nabla p^{n}=\boldsymbol{f}^{n} \tag{4.4}
\end{equation*}
$$

in $\Omega$ and that by (4.3) $\left\|\nabla p^{n}\right\|_{q} \leq c_{8}$. The sequence $\left\{\boldsymbol{u}^{n}\right\}$ is bounded in the space $D\left(A_{\gamma}^{\omega}\right)$. Hence there exists a subsequence again denoted by $\left\{\boldsymbol{u}^{n}\right\}$ which is weakly convergent in $D\left(A_{\gamma}^{\omega}\right)$. This subsequence naturally preserves the property (4.2).

Put $\boldsymbol{v}^{n}:=\left(\boldsymbol{u}^{n+1}-\boldsymbol{u}^{n}\right) / \delta_{n}$ where $\delta_{n}=\left\|\boldsymbol{u}^{n+1}-\boldsymbol{u}^{n}\right\|_{q}$. Then $\left\{\boldsymbol{v}^{n}\right\}$ is a sequence in the unit sphere in $L_{\sigma}^{q}(\Omega)$. It converges weakly to zero in $D\left(A_{\gamma}^{\omega}\right)$ because $\left(\boldsymbol{u}^{n+1}-\boldsymbol{u}^{n}\right) \rightharpoonup \mathbf{0}$ in $L_{\sigma}^{q}(\Omega)$ as $n \rightarrow \infty$ and by (4.1) $\delta_{n} \geq 1$. Hence $\left\{\boldsymbol{v}^{n}\right\}$ converges strongly to $\mathbf{0}$ in $W^{1, q}\left(\Omega \cap B_{R}(\mathbf{0})\right)^{3}$ for each $R>0$. Note that the function $\boldsymbol{v}^{n}$ satisfies the equation

$$
\begin{align*}
v \Delta \boldsymbol{v}^{n}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{v}^{n}-\boldsymbol{\omega} & \times \boldsymbol{v}^{n}-\gamma \partial_{1} \boldsymbol{v}^{n}-\lambda \boldsymbol{v}^{n}-\frac{1}{\delta_{n}} \nabla\left(p^{n+1}-p^{n}\right) \\
& =\frac{1}{\delta_{n}}\left(\boldsymbol{f}^{n+1}-\boldsymbol{f}^{n}\right) \tag{4.5}
\end{align*}
$$

in $\Omega$. This equation, together with the information on the weak convergence of $\left\{\boldsymbol{v}^{n}\right\}$ to zero in $D\left(A_{\gamma}^{\omega}\right)$, implies that the sequence $\left\{\nabla\left(p^{n+1}-p^{n}\right)\right\}$ weakly converges to zero in $L^{q}(\Omega)^{3}$. Thus, the functions $p^{n}$, which are given uniquely up to an additive constant, can be chosen so that $p^{n+1}-p^{n} \rightarrow 0$ strongly in $L^{q}\left(\Omega \cap B_{R}(\mathbf{0})\right)$ for each $R>0$.

The sequence $\left\{\boldsymbol{v}^{n}\right\}$ does not contain any subsequence, convergent in $L_{\sigma}^{q}(\Omega)$. Indeed, assume that $\left\{\boldsymbol{v}^{k_{n}}\right\}$ is a convergent subsequence of $\left\{\boldsymbol{v}^{n}\right\}$ in $L_{\sigma}^{q}(\Omega)$. This subsequence has the same weak limit as $\left\{\boldsymbol{v}^{n}\right\}$, hence $\boldsymbol{v}^{k_{n}} \rightharpoonup \mathbf{0}$ in $L_{\sigma}^{q}(\Omega)$ as $n \rightarrow \infty$. Then even $\boldsymbol{v}^{k_{n}} \rightarrow \mathbf{0}$ in $L_{\sigma}^{q}(\Omega)$ as $n \rightarrow \infty$. However, this is impossible because $\left\|\boldsymbol{v}^{k_{n}}\right\|_{q}=1$.

Suppose that $R>0$ is so large that the domain $\left\{\boldsymbol{x} \in \boldsymbol{R}^{3} ;|\boldsymbol{x}|>R\right\}$ is a subset of $\Omega$. Let $\eta$ be an infinitely differentiable cut-off function in $\Omega$, with values in [0, 1], such that

$$
\eta(\boldsymbol{x})= \begin{cases}0 & \text { for }|\boldsymbol{x}| \leq R, \\ 1 & \text { for } R+1 \leq|\boldsymbol{x}|\end{cases}
$$

Put $K_{R}:=\left\{\boldsymbol{x} \in \boldsymbol{R}^{3} ; R<|\boldsymbol{x}|<R+1\right\}$ and let $\mathfrak{B}: W_{0}^{1, q}\left(K_{R}\right) \mapsto W_{0}^{2, q}\left(K_{R}\right)^{3}$ be the Bogovskij operator (see the proof of Lemma 2.3). Then $\boldsymbol{V}^{n}:=\mathfrak{B}\left(\nabla \eta \cdot \boldsymbol{v}^{n}\right)$ belongs to $W_{0}^{2, q}\left(K_{R}\right)^{3}$. If we extend it by zero to $\Omega \backslash K_{R}$, it can be considered as an element of $W_{0}^{2, q}(\Omega)^{3}$. Due to the continuity of the operator $\mathfrak{B}$ and the strong convergence of $\left\{\boldsymbol{v}^{n}\right\}$ to $\mathbf{0}$ in $W^{1, q}\left(\Omega \cap B_{R+1}(\mathbf{0})\right)^{3}$, we get that $\boldsymbol{V}^{n} \rightarrow \mathbf{0}$ in $W_{0}^{2, q}(\Omega)^{3}$.

Now we define $\boldsymbol{w}^{n}(\boldsymbol{x}):=\eta(\boldsymbol{x}) \boldsymbol{v}^{n}(\boldsymbol{x})-\boldsymbol{V}^{n}(\boldsymbol{x})$. The function $\boldsymbol{w}^{n}$ belongs to $D\left(A_{\gamma}^{\omega}\right)$ and, due to (4.5), satisfies the equation

$$
\begin{align*}
& v \Delta \boldsymbol{w}^{n}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{w}^{n}-\boldsymbol{\omega} \times \boldsymbol{w}^{n}-\gamma \partial_{1} \boldsymbol{w}^{n}-\lambda \boldsymbol{w}^{n}-\frac{1}{\delta_{n}} \nabla\left[\eta\left(p^{n+1}-p^{n}\right)\right]  \tag{4.6}\\
& =\frac{\eta}{\delta_{n}}\left(\boldsymbol{f}^{n+1}-\boldsymbol{f}^{n}\right)-\frac{1}{\delta_{n}} \nabla \eta\left(p^{n+1}-p^{n}\right)+\nu(\Delta \eta) \boldsymbol{v}^{n}+2 v \nabla \eta \cdot \nabla \boldsymbol{v}^{n}-v \Delta \boldsymbol{V}^{n}+\lambda \boldsymbol{V}^{n} \\
& \quad+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot\left(\nabla \eta \otimes \boldsymbol{v}^{n}\right)-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{V}^{n}+\boldsymbol{\omega} \times \boldsymbol{V}^{n}-\gamma\left(\partial_{1} \eta\right) \boldsymbol{v}^{n}+\gamma \partial_{1} \boldsymbol{V}^{n} .
\end{align*}
$$

The right-hand side converges strongly to zero in $L^{q}(\Omega)^{3}$ as $n \rightarrow \infty$; this follows from the strong convergence of $\left\{\boldsymbol{v}^{n}\right\}$ to zero in $W^{1, q}\left(\Omega \cap B_{R+1}(\mathbf{0})\right)^{3}$, the strong convergence of $\left\{p^{n+1}-p^{n}\right\}$ to zero in $L^{q}\left(\Omega \cap B_{R+1}(\mathbf{0})\right)$, from the information on the support of $\nabla \eta$ and from the strong convergence of $\left\{\boldsymbol{V}^{n}\right\}$ to zero in $W^{2, q}(\Omega)^{3}$. Hence

$$
\begin{equation*}
\left\|\left(A_{\gamma}^{\omega}-\lambda I\right) \boldsymbol{w}^{n}\right\|_{q} \rightarrow 0 \quad \text { as } n \rightarrow \infty \tag{4.7}
\end{equation*}
$$

Moreover, given $\varepsilon>0$, there exists $n_{0} \in N$ such that if $n \in N, n \geq n_{0}$, then

$$
\begin{aligned}
\left\|\boldsymbol{w}^{n}\right\|_{q} & \leq\left(\int_{|x|<R+1}\left|\eta \boldsymbol{v}^{n}-\boldsymbol{V}^{n}\right|^{q} \mathrm{~d} \boldsymbol{x}\right)^{1 / q}+\left(\int_{R+1<|x|}\left|\boldsymbol{v}^{n}\right|^{q} \mathrm{~d} \boldsymbol{x}\right)^{1 / q} \leq \varepsilon+1 \\
\left\|\boldsymbol{w}^{n}\right\|_{q} & \geq\left(\int_{R+1<|x|}\left|\boldsymbol{v}^{n}\right|^{q} \mathrm{~d} \boldsymbol{x}\right)^{1 / q} \geq\left(\int_{\Omega}\left|\boldsymbol{v}^{n}\right|^{q} \mathrm{~d} \boldsymbol{x}\right)^{1 / q}-\left(\int_{|x|<R+1}\left|\boldsymbol{v}^{n}\right|^{q} \mathrm{~d} \boldsymbol{x}\right)^{1 / q} \\
& \geq 1-\varepsilon .
\end{aligned}
$$

Let us now normalize the sequence $\left\{\boldsymbol{w}^{n}\right\}$ by dividing each of the functions $\boldsymbol{w}^{n}$ by its norm in $L_{\sigma}^{q}(\Omega)$. In order to preserve a simple notation, we denote the normalized functions again by $\boldsymbol{w}^{n}$. If we finally put $\boldsymbol{w}^{n}(\boldsymbol{x})=\mathbf{0}$ for $\boldsymbol{x} \in \boldsymbol{R}^{3} \backslash \Omega$, we obtain a non-compact sequence in the unit sphere in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$, satisfying (4.7) with $\|\cdot\|_{q}$ being the norm in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. Let us denote, for a while, by $\left(A_{\gamma}^{\omega}\right)_{\boldsymbol{R}^{3}}$ the operator $A_{\gamma}^{\omega}$, considered in $L_{\sigma}^{q}\left(\boldsymbol{R}^{3}\right)$. The existence of the sequence
$\left\{\boldsymbol{w}^{n}\right\}$ with the above properties implies that nul ${ }^{\prime}\left(\left(A_{\gamma}^{\omega}\right)_{\boldsymbol{R}^{3}}-\lambda I\right)=\infty$. Hence $\lambda \notin \rho\left(\left(A_{\gamma}^{\omega}\right) \boldsymbol{R}^{3}\right)$, which, due to Theorem 1.1, yields $\lambda \in \Lambda_{\gamma}^{\omega}$.

Lemma 4.3. Let $\lambda \in C \backslash \Lambda_{\gamma}^{\omega}$. Then either $\lambda \in \sigma_{\mathrm{p}}\left(A_{\gamma}^{\omega}\right)$ for all $1<q<\infty$ or $\lambda \in \rho\left(A_{\gamma}^{\omega}\right)$ for all $1<q<\infty$. Moreover, $\lambda \in \rho\left(A_{\gamma}^{\omega}\right)$ when $\operatorname{Re} \lambda \geq 0$.

Proof. Assume that $1<q<\infty$ and that $\lambda \in C \backslash \Lambda_{\gamma}^{\omega}$ is an eigenvalue of $\left(A_{\gamma}^{\omega}\right)_{\Omega, q}$ with nonzero eigenfunction $\boldsymbol{v} \in D\left(\left(A_{\gamma}^{\omega}\right)_{\Omega, q}\right)$ ), where $\left(A_{\gamma}^{\omega}\right)_{\Omega, q}$ denotes the operator $A_{\gamma}^{\omega}$ on $L_{\sigma}^{q}(\Omega)$. Let $p$ with $\nabla p \in L^{q}(\Omega)^{3}$ be a corresponding pressure function. Using the cut-off function $\eta$, the set $K_{R}$ and Bogovskij's operator $\mathfrak{B}: W_{0}^{1, q}\left(K_{R}\right) \mapsto W_{0}^{2, q}\left(K_{R}\right)^{3}$ from the proof of Lemma 4.2, we get that $\boldsymbol{w}_{1}:=\eta \boldsymbol{v}-\boldsymbol{V}$ (where $\boldsymbol{V}=\mathfrak{B}(\nabla \eta \cdot \boldsymbol{v})$ in $K_{R}, \boldsymbol{V}=\mathbf{0}$ in $\boldsymbol{R}^{3} \backslash K_{R}$ ) solves the equation

$$
\begin{align*}
& v \Delta \boldsymbol{w}_{1}+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{w}_{1}-\boldsymbol{\omega} \times \boldsymbol{w}_{1}-\gamma \partial_{1} \boldsymbol{w}_{1}-\lambda \boldsymbol{w}_{1}-\nabla(\eta p) \\
& =\boldsymbol{f}:=-\nabla \eta p+\nu(\Delta \eta) \boldsymbol{v}+2 v \nabla \eta \cdot \nabla \boldsymbol{v}-v \Delta \boldsymbol{V}+\lambda \boldsymbol{V}  \tag{4.8}\\
& \quad+(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot(\nabla \eta \otimes \boldsymbol{v})-(\boldsymbol{\omega} \times \boldsymbol{x}) \cdot \nabla \boldsymbol{V}+\boldsymbol{\omega} \times \boldsymbol{V}-\gamma\left(\partial_{1} \eta\right) \boldsymbol{v}+\gamma \partial_{1} \boldsymbol{V}
\end{align*}
$$

in $\boldsymbol{R}^{3}$ (cf. (4.6)).
First assume that $q \geq 3$. Due to the boundedness of the operator $\mathfrak{B}$ from $W_{0}^{2, q}\left(K_{R}\right)$ to $W_{0}^{3, q}\left(K_{R}\right)^{3}$ (see e.g. [15, p. 130]), the restriction of the vector field $\boldsymbol{V}$ to $K_{R}$ belongs to $W_{0}^{3, q}\left(K_{R}\right)^{3}$. From this, we deduce that $f \in W^{1, q}\left(K_{R}\right)$. Hence, by Sobolev's embedding theorem, $\boldsymbol{f} \in L^{s}\left(K_{R}\right)^{3}$ for $1<s<\infty$. Since $\boldsymbol{f}$ is supported in $\bar{K}_{R}$, we have that $\boldsymbol{f} \in$ $L^{s}\left(\boldsymbol{R}^{3}\right)^{3}$ for $1<s<\infty$. Now we may apply Theorem 1.1 (with $s$ instead of $q$ ) to the whole space problem (4.8) and conclude that $\boldsymbol{w}_{1} \in D\left(\left(A_{\gamma}^{\omega}\right)_{\boldsymbol{R}^{3}, s}\right) \subset W^{2, s}\left(\boldsymbol{R}^{3}\right)$ and $\nabla(\eta p) \in L^{s}\left(\boldsymbol{R}^{3}\right)^{3}$ because $\lambda \in \boldsymbol{C} \backslash \Lambda_{\gamma}^{\omega}$ belongs to the resolvent set of $\left(A_{\gamma}^{\omega}\right)_{\boldsymbol{R}^{3}, s}$. Similarly, we derive that $\boldsymbol{w}_{2}:=$ $(1-\eta) \boldsymbol{v}+\boldsymbol{V}$, the solution of a problem analogous to (4.8) in the bounded domain $\Omega_{R+1}:=$ $\Omega \cap B_{R+1}(\mathbf{0})$, satisfies $\boldsymbol{w}_{2} \in W^{2, s}\left(\Omega_{R+1}\right)$. Consequently, $\boldsymbol{v} \equiv \boldsymbol{w}_{1}+\boldsymbol{w}_{2} \in D\left(\left(A_{\gamma}^{\omega}\right)_{\Omega, s}\right)$ and $\lambda$ is an eigenvalue of $\left(A_{\gamma}^{\omega}\right)_{\Omega, s}$.

If $1<q<3$ then we obtain the same result for $1<s<3 q /(3-q)$. However, repeating finitely many times the same argument, we can extend the result to all $1<s<\infty$.

Finally, when $s=2$, a variational argument implies that $\operatorname{Re} \lambda<0$ for all $\lambda \in \sigma\left(A_{\gamma}^{\omega}\right)$ (cf. [14, Theorem 1.1]).

Now Theorem 1.2 is completely proved.
REMARK 4.1. If $q=2$ then the interesting question occurs whether $A_{\gamma}^{\omega}$ is a normal operator in $L_{\sigma}^{2}(\Omega)$. We have proved in our previous papers [13] and [14] that
(a) if $\gamma=0$ and the domain $\Omega$ is axially symmetric with respect to the $x_{1}$-axis then $A_{\gamma}^{\omega}$ is normal,
(b) if $\gamma \neq 0$ or the domain $\Omega$ is not axially symmetric then $A_{\gamma}^{\omega}$ is not normal.

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## References

[1] W. Borchers and H. Sohr, On the equations rot $v=g$ and $\operatorname{div} u=f$ with zero boundary conditions, Hokkaido Math. J. 19 (1990), 67-87.
[2] M. E. BogovskiJ, Solution of the first boundary value problem for an equation of continuity of an incompressible medium, Dokl. Akad. Nauk SSSR 248 (1979), 1037-1040.
[3] R. FARWIG and H. Sohr, Generalized resolvent estimates for the Stokes system in bounded and unbounded domains, J. Math. Soc. Japan 46 (1994), 608-643.
[4] R. FARWIG, An $L^{p}$-analysis of viscous fluid flow past a rotating obstacle, Tohoku Math. J. 58 (2005), 129147.
[5] R. FARWIG, Estimates of lower order derivatives of viscous fluid flow past a rotating obstacle, 73-84, Banach Center Publ. 70, Warsaw, 2005.
[6] R. FARWIG and T. Hishida, Asymptotic profile of steady Stokes flow around a rotating obstacle, Dept. of Math., Darmstadt University of Technology, preprint 2578 (2009), Manuscripta Math. (to appear).
[7] R. Farwig and T. Hishida, Asymptotic profiles of steady Stokes and Navier-Stokes flows around a rotating obstacle, Ann. Univ. Ferrara 55 (2009), 263-277.
[8] R. Farwig, T. Hishida and D. Müller, $L^{q}$-theory of a singular "winding" integral operator arising from fluid dynamics, Pacific J. Math. 215 (2004), 297-312.
[9] R. FARWIG, M. Krbec and Š. NeČasová, A weighted $L^{q}$-approach to Stokes flow around a rotating body, Ann. Univ. Ferrara Sez. VII Sci. Mat. 54 (2008), 61-84.
[10] R. Farwig, M. Krbec and Š. NečASová, A weighted $L^{q}$-approach to Oseen flow past a rotating body, Math. Methods Appl. Sci. 31 (2008), 551-574.
[11] R. FARWIG, Š. NEČASOVÁ AND J. NEUSTUPA, On the essential spectrum of a Stokes-type operator arising from flow around a rotating body in the $L^{q}$-framework, Kyoto Conference on the Navier-Stokes Equations and their Applications 2006, 93-105, RIMS Kôkyûroku Bessatsu B1, Res. Inst. Math. Sci. (RIMS), Kyoto, 2007.
[12] R. FARWIG, Š. NečASOVÁ and J. Neustupa, Spectral analysis of a Stokes-type operator from flow around a rotating body, Dept. of Math., Darmstadt University of Technology, preprint 2598 (2009).
[13] R. FARWIG and J. Neustupa, On the spectrum of a Stokes-type operator arising from flow around a rotating body, Manuscripta Math. 122 (2007), 419-437.
[14] R. FARWIG AND J. NEUSTUPA, On the spectrum of an Oseen-type operator arising from flow around a rotating body, Integral Equations and Operator Theory 62 (2008), 169-189.
[15] G. P. GALDI, An introduction to the mathematical theory of the Navier-Stokes equations, Vol. I, Linear steady problems, Springer Tracts in Natural Philosophy 38, Springer-Verlag, New York, 1994.
[16] G. P. Galdi, On the motion of a rigid body in a viscous liquid: A mathematical analysis with applications (S. Friedlander, D. Serre eds), Handbook of mathematical fluid dynamics, Vol. I, 653-791, North-Holland, Amsterdam, 2002.
[17] G. P. Galdi, Steady flow of a Navier-Stokes fluid around a rotating obstacle, J. Elasticity 71 (2003), 1-32.
[18] G. P. Galdi and M. Kyed, Steady-state Navier-Stokes flows past a rotating body: Leray solutions are physically reasonable, Institut für Mathematik, RWTH Aachen, preprint (2010).
[19] G. P. Galdi and M. Kyed, Asymptotic behavior of a Leray solution around a rotating obstacle, Institut für Mathematik, RWTH Aachen, preprint (2010).
[20] M. Geissert, H. Heck and M. Hieber, $L^{p}$-theory of the Navier-Stokes flow in the exterior of a moving or rotating obstacle, J. Reine Angew. Math. 596 (2006), 45-62.
[21] T. Hishida, An existence theorem for the Navier-Stokes flow in the exterior of a rotating obstacle, Arch. Rational Mech. Anal. 150 (1999), 307-348.
[22] T. Hishida, The Stokes operator with rotating effect in exterior domains, Analysis 19 (1999), 51-67.
[23] T. Hishida, $L^{q}$ estimates of weak solutions to the stationary Stokes equations around a rotating body, J. Math. Soc. Japan 58 (2006), 743-767.
[24] T. Hishida and Y. Shibata, $L_{p}-L_{q}$ estimate of the Stokes operator and Navier-Stokes flows in the exterior of a rotating obstacle, Arch. Ration. Mech. Anal. 193 (2009), 339-421.
[25] T. Kato, Perturbation theory for linear operators, Springer-Verlag, Berlin-Heidelberg-New York, 1966.
[26] A. Kufner, S. Fučík and O. John, Function spaces, Academia, Praha, 1977.
[27] O. A. LADYZhenskaya, The mathematical theory of viscous incompressible flow, Gordon and Breach, Science Publishers, New York, 1969.
[28] Y. Shibata, On the Oseen semigroup with rotating effect, Functional analysis and evolution equations, The Günter Lumer Volume, (Eds. H. Amann, W. Arendt, M. Hieber, F. Neubrander, B. Nicaise, J. von Below), 595-611, Birkhäuser, Basel, 2008.

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