# $L^{p}$-spectrum of degenerate hypoelliptic Ornstein-Uhlenbeck operators 

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#### Abstract

We describe the spectrum of degenerate hypoelliptic Ornstein-Uhlenbeck operators $\mathcal{A}=\sum_{i, j=1}^{n} q_{i j} D_{i j}+\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}$ in $L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<+\infty$, and in $C_{0}\left(\mathbb{R}^{n}\right)$. We show that the spectrum of $\mathcal{A}$ is the sum of $(-\infty, 0]$ and the spectrum of the drift term. Our result gives a complete picture of the spectral properties of Ornstein-Uhlenbeck operators in $L^{p}$ spaces.


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## 1 Introduction

The aim of the present paper is the full description of the spectrum of possibly degenerate Ornstein-Uhlenbeck operators

$$
\begin{equation*}
\mathcal{A}=\sum_{i, j=1}^{n} q_{i j} D_{i j}+\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}=\operatorname{Tr}\left(Q D^{2}\right)+\langle B x, D\rangle, \quad x \in \mathbb{R}^{n} \tag{1.1}
\end{equation*}
$$

in $L^{p}\left(\mathbb{R}^{n}\right), 1 \leq p<+\infty$, and in $C_{0}\left(\mathbb{R}^{n}\right)$. Here $Q=\left(q_{i j}\right)$ is a real, constant, symmetric and positive semidefinite matrix and $B=\left(b_{i j}\right)$ is a nonzero real matrix. The semidefinitess of the matrix $Q$ is responsible for the possible degeneracy of $\mathcal{A}$. Throughout we assume that $\mathcal{A}$ is hypoelliptic, which can be stated as follows: the symmetric matrices

$$
Q_{t}=\int_{0}^{t} e^{s B} Q e^{s B^{T}} d s
$$

have nonzero determinant for some (equivalently, for all) $t>0$. In the literature one can find several equivalent conditions for hypoellipticity. In particular, on p. 148 of [10] it is pointed out that the hypoellipticity of $\mathcal{A}$ is equivalent to the property

$$
\begin{equation*}
\operatorname{ker}(Q) \text { does not contain nontrivial subspaces which are invariant for } B^{T} \tag{1.2}
\end{equation*}
$$

see also [13, Appendix]. Here nontrivial means different from $\{0\}$.

[^0]The hypoellipticity assumption implies that the associated Markov semigroup $(T(t))_{t \geq 0}$ has the following explicit representation formula due to Kolmogorov [12]

$$
\begin{equation*}
(T(t) f)(x)=\frac{1}{(4 \pi)^{n / 2}\left(\operatorname{det} Q_{t}\right)^{1 / 2}} \int_{\mathbb{R}^{n}} e^{-\left\langle Q_{t}^{-1} y, y\right\rangle / 4} f\left(e^{t B} x-y\right) d y, \quad x \in \mathbb{R}^{n}, t>0 \tag{1.3}
\end{equation*}
$$

The parabolic equation $u_{t}=\mathcal{A} u$, known as Kolmogorov equation, is solved by the function $u(t, x)=T(t) f(x)$ for a large class of initial data $f$. In recent years, both the semigroup $(T(t))_{t \geq 0}$ and its generator $\mathcal{A}$ have extensively been studied. Several applications in physics and finance for the operator $\mathcal{A}$ and its evolutionary counterpart $\mathcal{A}-D_{t}$ can be found in the survey [21]. They were also used in the context of rotating fluids, see e.g. [9]. These operators were also the leading example for an intensive research on elliptic and parabolic problems with unbounded coefficients, see e.g. [14].

In the analytical study of $\mathcal{A}$, even in the nondegenerate case the classical $L^{p}$ and Schauder estimates do not apply because of the unboundedness of the first order coefficients. Regularity properties in spaces of continuous functions were proved in [4] in the nondegenerate case and in [15] in the degenerate case. Schauder estimates can then be deduced by means of interpolation techniques. Moreover, $L^{p}$ estimates were established in [20] and in [19] in the nondegenerate case, by a semigroup approach, and in [1] in the degenerate case.

The underlying stochastic process admits an invariant measure $\mu$ if and only if all eigenvalues of the drift matrix $B$ have negative real parts. This means that $\mu$ is a probability measure satisfying

$$
\int_{\mathbb{R}^{n}}(T(t) f)(x) d \mu(x)=\int_{\mathbb{R}^{n}} f(x) d \mu(x)
$$

for every $t \geq 0$ and continuous and bounded function $f$ on $\mathbb{R}^{n}$. The invariant measure is unique and absolutely continuous with respect to the Lebesgue measure having the (Gaussian) density

$$
\rho(x)=\frac{1}{(4 \pi)^{n / 2}\left(\operatorname{det} Q_{\infty}\right)^{1 / 2}} e^{-\left\langle Q_{\infty}^{-1} x, x\right\rangle / 4} \quad \text { with } \quad Q_{\infty}=\int_{0}^{\infty} e^{s B} Q e^{s B^{T}} d s
$$

see [5, Chapter II. 6].
The semigroup $(T(t))_{t \geq 0}$ and its generator $\mathcal{A}$ have widely been investigated in the weighted spaces $L^{p}\left(\mathbb{R}^{n}, d \mu\right)$, if $\sigma(B) \subset \mathbb{C}_{-}$. Here the unboundedness of the coefficients of $\mathcal{A}$ is balanced by the exponential decay of the density $\rho$ which leads to a much better behavior in several respects. For instance, the generator has compact resolvent in $L^{p}\left(\mathbb{R}^{n}, d \mu\right)$ if $p \in(1, \infty)$, which is not true in the unweighted spaces $L^{p}$. The domain of the generator in $L^{p}\left(\mathbb{R}^{n}, d \mu\right)$ was computed in [16] for $p=2$ and in [20] for $p \in(1, \infty)$ in the nondegenerate case. See also $[2,3]$ for the analogous problem on an infinite-dimensional Hilbert space $E$ instead of $\mathbb{R}^{n}$. In the degenerate case a sharp inclusion for the domain was shown in [8] for $p=2$, whereas the case $p \neq 2$ is still an open problem, indicating that the general picture of Ornstein-Uhlenbeck operators is still not complete.

In [18] the spectrum of $\mathcal{A}$ in $L^{p}\left(\mathbb{R}^{n}, d \mu\right)$ was completely described also in the degenerate case, provided that $\sigma(B) \subset \mathbb{C}_{-}$. The situation is much different in the spaces $L^{p}=L^{p}\left(\mathbb{R}^{n}\right)$ with respect to the Lebesgue measure, e.g., since $\mathcal{A}$ does not have a compact resolvent here. For some choices of $B$ the spectrum of $\mathcal{A}$ was computed in $L^{p}$ in [17]. This paper is the starting point of our investigation.

The operator $\mathcal{A}$ can be seen as the sum of the diffusion term $\sum_{i, j=1}^{n} q_{i j} D_{i j}$ and of the drift $\operatorname{term} \mathcal{L}=\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}$. The spectral properties of the drift term are fully understood
by [17]. There it was proved that the spectrum of the realization $\mathcal{L}_{p}$ of $\mathcal{L}$ in $L^{p}$ is the line $-\operatorname{tr}(B) / p+i \mathbb{R}$ unless $B$ is similar to a diagonal matrix with purely imaginary eigenvalues. In this last case the spectrum of $\mathcal{L}_{p}$ can be either $i \mathbb{R}$ or a discrete, explicitly given subgroup $G$ of $i \mathbb{R}$, see Theorem 2.2 and Proposition 2.3.

In [17] it is further shown that the boundary spectrum of the realization $\mathcal{A}_{p}$ of $\mathcal{A}$ in $L^{p}$ contains the spectrum of $\mathcal{L}_{p}$ without further assumptions on the matrices $Q$ and $B \neq 0$. Here $\mathcal{A}_{p}$ is defined as the generator of $(T(t))_{t \geq 0}$ in $L^{p}$, see Proposition 2.4. The spectrum of $\mathcal{A}_{p}$ has been computed in [17] if $\sigma(B)$ is contained in the left or in the right open half-plane. In this case, $\sigma\left(\mathcal{A}_{p}\right)$ is equal to $\{\mu \in \mathbb{C}: \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$. So it depends on $p$ and is far from being discrete. In addition, and this is the main step in [17], if all the eigenvalues of $B$ have positive real parts, then the open half-plane $\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu<-\operatorname{tr}(B) / p\}$ consists of eigenvalues.

In this paper we complete the picture computing the spectrum of $\mathcal{A}_{p}$ without any further restriction on $Q=Q^{T} \geq 0$ and $B \neq 0$, apart from hypoellipticity. We prove that $\sigma\left(\mathcal{A}_{p}\right)$ is given as the sum of the spectra of its diffusion part (i.e., $(-\infty, 0])$ and of the drift term $\mathcal{L}_{p}$.
Theorem 1.1. Let (1.2) be true and $p \in[1, \infty]$. Then the spectrum of $\mathcal{A}_{p}$ is given by

$$
\sigma\left(\mathcal{A}_{p}\right)=(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)
$$

In particular, we have either $\sigma\left(\mathcal{A}_{p}\right)=(-\infty, 0]+G$ or $\sigma\left(\mathcal{A}_{p}\right)=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$, according to $\sigma\left(\mathcal{L}_{p}\right)$ being a discrete subgroup $G=\frac{2 \pi i}{\tau} \mathbb{Z}$ of $i \mathbb{R}$ or the whole line $-\operatorname{tr}(B) / p+i \mathbb{R}$. Moreover, the semigroup $(T(t))_{t \geq 0}$ satisfies the weak spectral mapping theorem

$$
\sigma(T(t))=\overline{\exp \left(t \sigma\left(\mathcal{A}_{p}\right)\right)}, \quad t \geq 0
$$

We even have $\sigma(T(t)) \backslash\{0\}=e^{t \sigma\left(\mathcal{A}_{p}\right)}$ except for the case that $\sigma\left(\mathcal{L}_{p}\right)=G=\frac{2 \pi i}{\tau} \mathbb{Z}$ and $t / \tau$ is irrational.

We note that for $p=2$ the spectral mapping theorem was proved for perturbed OrnsteinUhlenbeck operators with $Q=I$ and $B=2 I$ by completely different methods in [6].

If $B=B^{T}$ and $Q B=B Q$, by separation of variables one can transform the OrnsteinUhlenbeck operator into the form $\mathcal{A}=\Delta+\sum_{i=1}^{n} b_{i} x_{i} D_{i}$. Here the problem can be reduced to one dimensional problems, see [17, Theorem 5.1]. However, this is far from being the general case. We also stress that $\mathcal{A}_{2}$ does not possess eigenvalues if $B$ has an eigenvalue with negative real part or if $B$ is skew-symmetric and $Q=I$, as we will see in Section 3 . So we have to proceed in a different way than in [17] or [18], where eigenfunctions played a crucial role.

Instead, we start by reducing $\mathcal{A}$ to a canonical form with an upper quasi triangular drift matrix whose diagonal is formed by $1 \times 1$ and $2 \times 2$ blocks containing the real and complex conjugate eigenvalues of $B$, respectively. The transformation is made through a linear change of variables that leaves the spectrum unchanged.

The second step consists in a scaling procedure leading to a new operator $\mathcal{C}$ in the limit which is the sum of an Ornstein-Uhlenbeck operator in one or two variables and a drift operator acting in the remaining ones. The scaling and the limit allow us to get rid of the upper off-diagonal blocks of the drift matrix of $\mathcal{A}$ and to separate the variables. We can recover the spectrum of $\mathcal{A}_{p}$ from that of the limit operator $\mathcal{C}$.

The main part of the proof is thus devoted to the investigation of the spectrum of $\mathcal{C}$. Here we can assume that $B$ has an eigenvalue with nonnegative real part, since the other case is already covered by the main result in [17]. The above splitting then reduces the problem
to Ornstein-Uhlenbeck operators in $\mathbb{R}$ or in $\mathbb{R}^{2}$ where $B$ has one nonnegative eigenvalue or two complex conjugate eigenvalues with nonnegative real parts. We further have to treat eigenvalues in $i \mathbb{R}$ and with positive real part separately. The detailed study of these four cases is mainly based on the construction of approximate eigenfuctions.

The paper is structured as follows. In Section 2 we recall the known generator properties of the drift operator $\mathcal{L}=\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}$ and its spectrum, as computed in [17]. We provide further details in the case where the generated group is periodic. We also collect the known properties on $\mathcal{A}$. Most of the results are contained in [17], where it is assumed that $Q$ is positive definite. However, we explain why they continue to hold with minor modifications in the degenerate hypoelliptic setting. Corollary 2.7 and Proposition 2.8 establish the inclusion $\sigma\left(\mathcal{A}_{p}\right) \subseteq(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$ by means of general spectral theory of semigroups. In Section 3 we show that there are no eigenvalues of $\mathcal{A}_{2}$ in many cases. Finally, Section 4 is devoted to the proof of Theorem 1.1. Here also the spectral mapping theorem follows mainly from general theory, whereas the proof of the other inclusion $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$ requires a sophisticated analysis of the four cases indicated above.

Warning: Throughout the whole paper, we write $L^{\infty}$ for $C_{0}\left(\mathbb{R}^{n}\right)$, which is the space of continuous functions on $\mathbb{R}^{n}$ vanishing at infinity, endowed with the supremum norm.

Notation. $L^{p}$ stands for $L^{p}\left(\mathbb{R}^{n}\right)$ if $p \in[1, \infty)$ and $C_{c}^{\infty}$ for $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$.
The spectrum and the resolvent set of a linear operator $\mathcal{B}$ are denoted by $\sigma(\mathcal{B})$ and $\rho(\mathcal{B})$, respectively. The spectral bound of $\mathcal{B}$ is defined by $s(\mathcal{B})=\sup \{\operatorname{Re} \mu \mid \mu \in \sigma(\mathcal{B})\}$ and the boundary spectrum is $\sigma(\mathcal{B}) \cap\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu=s(\mathcal{B})\}$. The approximate point spectrum $\sigma_{\text {ap }}(\mathcal{B})$ of $\mathcal{B}$ is the subset of $\sigma(\mathcal{B})$ of all complex numbers $\mu$ for which there is a sequence $\left(v_{n}\right)$ in its domain $D(\mathcal{B})$ such that $\left\|v_{n}\right\|=1$ and $\left\|\mathcal{B} v_{n}-\mu v_{n}\right\| \rightarrow 0$ as $n \rightarrow \infty$. The sequence $\left(v_{n}\right)$ is called an approximate eigenvector relative to the approximate eigenvalue $\mu$. The topological boundary of the spectrum of $\mathcal{B}$ is always contained in $\sigma_{a p}(\mathcal{B})$ (see [7, Proposition IV.1.10]).

We write $\mathcal{B}_{p}$ to indicate a realization of a (differential) operator $\mathcal{B}$ in $L^{p}$, that is when $\mathcal{B}$ is provided with a specific domain in $L^{p}$. However, we sometimes omit the suffix $p$ in the proofs, to shorten the notation.

If $B$ is a matrix, $B^{T}$ denotes its transpose. We set $\mathbb{C}_{+}=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu>0\}$ and $\mathbb{C}_{-}=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu<0\}$. When $p=\infty$, by $1 / p$ we mean 0 .

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## 2 Basic and known results

We collect background material from [17] and prove auxiliary results concerning the drift term and the Ornstein-Uhlenbeck operator.

### 2.1 Properties of $\mathcal{L}$

Let $B=\left(b_{i j}\right) \neq 0$ be a real $n \times n$ matrix and consider the drift operator

$$
\mathcal{L}=\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}
$$

defined on its maximal domain

$$
D\left(\mathcal{L}_{p}\right)=\left\{u \in L^{p} \mid \mathcal{L} u \in L^{p}\right\}
$$

in $L^{p}$ for $1 \leq p \leq \infty$, where $\mathcal{L} u$ is understood in the sense of distributions. We write $\mathcal{L}_{p}$ for $\left(\mathcal{L}, D\left(\mathcal{L}_{p}\right)\right)$ and recall the following results, whose proofs can be found in [17, Section 2$]$.

Proposition 2.1. Let $1 \leq p \leq \infty$. The operator $\mathcal{L}_{p}$ generates the $C_{0}$-group $(S(t))_{t \in \mathbb{R}}$ in $L^{p}$ defined by

$$
\begin{equation*}
(S(t) f)(x)=f\left(e^{t B} x\right) \tag{2.1}
\end{equation*}
$$

the space $C_{c}^{\infty}$ is a core of $\mathcal{L}_{p}$, and we have

$$
\begin{equation*}
\|S(t) f\|_{p}=e^{-\frac{t}{p} \operatorname{tr}(B)}\|f\|_{p} \tag{2.2}
\end{equation*}
$$

for $f \in L^{p}$ and $t \in \mathbb{R}$.
We next describe the spectrum of $\mathcal{L}_{p}$ distinguishing several cases.
Theorem 2.2. Let $1 \leq p \leq \infty$.
(a) Let $\operatorname{tr}(B) \neq 0$. Then $\sigma\left(\mathcal{L}_{p}\right)=-\operatorname{tr}(B) / p+i \mathbb{R}$.
(b) Let $\operatorname{tr}(B)=0$ and $B$ be not similar to a diagonal matrix with purely imaginary eigenvalues. Then $\sigma\left(\mathcal{L}_{p}\right)=i \mathbb{R}$.
(c) Let $B$ be similar to a diagonal matrix with nonzero eigenvalues $\pm i \sigma_{1}, \pm i \sigma_{2}, \ldots, \pm i \sigma_{k}$ in $i \mathbb{R}$ and possibly 0 , where $\sigma_{r} \sigma_{s}^{-1} \notin \mathbb{Q}$ for some $r, s \in\{1, \ldots, k\}$. Then $\sigma\left(\mathcal{L}_{p}\right)=i \mathbb{R}$.
(d) Let $B$ be similar to a diagonal matrix with nonzero eigenvalues $\pm i \sigma_{1}, \pm i \sigma_{2}, \ldots, \pm i \sigma_{k}$ in $i \mathbb{R}$ and possibly 0 , where $\sigma_{r} \sigma_{s}^{-1} \in \mathbb{Q}$ for all $r, s \in\{1, \ldots, k\}$. Then $(S(t))_{t \in \mathbb{R}}$ is periodic and $\sigma\left(\mathcal{L}_{p}\right)$ is the discrete subgroup $G=\left\{i\left(n_{1} \sigma_{1}+\cdots+n_{k} \sigma_{k}\right) \mid\left(n_{1}, \ldots, n_{k}\right) \in \mathbb{Z}^{k}\right\}$.

In the sequel we need more information about case (d) above in which $(S(t))_{t \in \mathbb{R}}$ is periodic.
Proposition 2.3. Let $B$ be similar to a diagonal matrix with nonzero eigenvalues $\pm i \sigma_{1}, \pm i \sigma_{2}, \ldots, \pm i \sigma_{k}$ in $i \mathbb{R}$ and possibly 0 , with $2 k \leq n$. Assume that for every $j \in$ $\{2, \ldots, k\}$ we have $\sigma_{j}=\frac{p_{j}}{q_{j}} \sigma_{1}$ for some coprime integers $p_{j}$ and $q_{j}$. Then $(S(t))$ is periodic with period $\tau=2 \pi N \sigma_{1}^{-1}$, where $N$ is the least common multiple of $q_{2}, \ldots, q_{k}$. Moreover, the set $G$ from Theorem 2.2 is given by $G=\frac{\sigma_{1}}{N} i \mathbb{Z}=\frac{2 \pi}{\tau} i \mathbb{Z}$.

Proof. We denote a point in $\mathbb{R}^{n}$ by $x=\left(x_{1}, y_{1}, \ldots, x_{k}, y_{k}, w_{2 k+1}, \ldots, w_{n}\right)$ and set $z_{j}=$ $\left(x_{j}, y_{j}\right)$. Possibly after a change of variables we obtain

$$
\begin{equation*}
S(t) f(x)=f\left(e^{i t \sigma_{1}} z_{1}, \ldots, e^{i t \sigma_{k}} z_{k}, w_{2 k+1}, \ldots, w_{n}\right) \tag{2.3}
\end{equation*}
$$

see Theorem 2.6 of [17]. If $0 \notin \sigma(B)$, the components $w_{j}$ are not present. Formula (2.3) yields $S(\tau) f=f$.

We prove that the set $G$ defined in Theorem $2.2(\mathrm{~d})$ coincides with $\frac{\sigma_{1}}{N} i \mathbb{Z}$. The inclusion $\subseteq$ easily follows from the form of the numbers $\sigma_{j}$ and the definition of $N$. To show the other inclusion, we first observe that the greatest common divisor of $N, N p_{2} / q_{2}, \ldots, N p_{k} / q_{k}$ is equal to 1 . Indeed, otherwise there would exist a prime number $p$ dividing $N, \ldots, N p_{k} / q_{k}$. Let $\alpha \in \mathbb{N}$ be the greatest exponent for which $p^{\alpha}$ divides $N$. Then $p^{\alpha}$ occurs in the prime
factorization of some $q_{j}$. Since $p_{j}$ and $q_{j}$ are coprime, $p$ cannot divide $N p_{j} / q_{j}$, and this is a contradiction. As a result, each integer $m$ can be written as

$$
m=m_{1} N+m_{2} \frac{N p_{2}}{q_{2}}+\cdots+m_{k} \frac{N p_{k}}{q_{k}}
$$

for suitable $m_{j} \in \mathbb{Z}$. This is equivalent to saying that the element $\frac{\sigma_{1}}{N} m$ can be written as $m_{1} \sigma_{1}+\cdots+m_{k} \sigma_{k}$ and concludes the proof.

### 2.2 Properties of $\mathcal{A}$

We turn our attention to the Ornstein-Uhlenbeck operator defined in (1.1) and to the associated semigroup $(T(t))_{t \geq 0}$ given by (1.3). We always assume the hypoellipticity condition (1.2) and $1 \leq p \leq \infty$. We do not need the full description of the domain of the generator, but only the fact that smooth functions with compact support are a core. We point out, however, that the domain has been described in [20, Section 4] and in [19] in the nondegenerate case and in [1] in the degenerate one.

Proposition 2.4. The semigroup $(T(t))_{t \geq 0}$ is strongly continuous on $L^{p}, 1 \leq p \leq \infty$, and satisfies the estimate

$$
\begin{equation*}
\|T(t)\| \leq e^{-\frac{t}{p} \operatorname{tr}(B)} \tag{2.4}
\end{equation*}
$$

for $t \geq 0$. Moreover, $C_{c}^{\infty}$ is a core for the generator $\mathcal{A}_{p}$.
Proof. If the diffusion matrix $Q$ is positive definite, the stated properties and a partial description of the domain of the generator have been proved in Section 3 of [17]. However, the same proofs hold in the degenerate hypoelliptic case. We only sketch them and refer to [17] for more details. To show (2.4), we write $T(t) f=S(t)\left(g_{t} * f\right)$ where

$$
g_{t}(y)=\frac{1}{(4 \pi)^{n / 2}\left(\operatorname{det} Q_{t}\right)^{1 / 2}} e^{-\left\langle Q_{t}^{-1} y, y\right\rangle / 4}
$$

and $S(t)$ is defined in (2.1). The estimate (2.4) then follows from (2.2), Young's inequality for convolutions, and $\left\|g_{t}\right\|_{1}=1$. Since $T(t) f \rightarrow f$ in $L^{p}$ as $t \rightarrow 0^{+}$for $f \in C_{c}^{\infty}$, by density (2.4) implies the strong continuity of $(T(t))_{t \geq 0}$ for every $1 \leq p \leq \infty$.

Let $\mathcal{A}_{p}$ and $\mathcal{S}\left(\mathbb{R}^{n}\right)$ denote the generator of $(T(t))_{t \geq 0}$ in $L^{p}$ and the Schwartz class, respectively. One easily checks that $\mathcal{S}\left(\mathbb{R}^{n}\right) \subseteq D\left(\mathcal{A}_{p}\right)$ and $\mathcal{A}_{p} f=\mathcal{A} f$ for every $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$. Moreover, $\mathcal{S}\left(\mathbb{R}^{n}\right)$ is dense in $L^{p}$ and invariant for $(T(t))_{t \geq 0}$ by (1.3). Therefore it is a core of $\mathcal{A}_{p}$. By a truncation argument we conclude that $C_{c}^{\infty}$ is a core for $\mathcal{A}_{p}$.

We recall Theorem 3.3 and Corollary 3.5 of [17].
Proposition 2.5. The boundary spectrum of $\mathcal{A}_{p}$ contains the spectrum of the drift $\mathcal{L}_{p}$.
Corollary 2.6. The growth bound of $(T(t))_{t \geq 0}$ in $L^{p}$ is $\omega_{p}=-\operatorname{tr}(B) / p$.
Standard semigroup theory then yields first inclusions of the spectra.
Corollary 2.7. The spectrum of $\mathcal{A}_{p}$ belongs to the half-plane $\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$, and that of $T(t)$ to the closed ball $\bar{B}\left(0, e^{-\frac{t}{p} \operatorname{tr}(B)}\right)$.

If $\sigma\left(\mathcal{L}_{p}\right)$ is the whole line $-\operatorname{tr}(B) / p+i \mathbb{R}$, the half-plane $\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$ coincides with the sum $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$. This is not the case if $\sigma\left(\mathcal{L}_{p}\right)$ is a discrete subgroup of $i \mathbb{R}$, which occurs when the group generated by the drift part is periodic (see Theorem $2.2(\mathrm{~d}))$. However, also in this case we have the inclusion $\sigma\left(\mathcal{A}_{p}\right) \subseteq(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$ as proved in the next proposition.

Proposition 2.8. Let $B$ be similar to a diagonal matrix with nonzero eigenvalues $\pm i \sigma_{1}$, $\cdots, \pm i \sigma_{k}$ in $i \mathbb{R}$ and possibly 0 . Assume that the quotient $\sigma_{r} \sigma_{s}^{-1}$ is rational for all $r$ and $s$. Then $\sigma\left(\mathcal{A}_{p}\right) \subseteq(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$.
Proof. Let $\tau>0$ be the period of $(S(t))_{t \in \mathbb{R}}$, see Proposition 2.3. Also $e^{t B^{T}}$ is $\tau$-periodic and in particular $e^{\tau B^{T}}=I$. By the representation formula (1.3) we have

$$
(T(\tau) f)(x)=\frac{1}{(4 \pi)^{n / 2}\left(\operatorname{det} Q_{\tau}\right)^{1 / 2}} \int_{\mathbb{R}^{n}} e^{-\left\langle Q_{\tau}^{-1} y, y\right\rangle / 4} f(x-y) d y
$$

showing that $T(\tau)=T_{\tau}(1)$ where $\left(T_{\tau}(t)\right)_{t \geq 0}$ is the semigroup generated by the diffusion operator $A_{\tau}=\operatorname{Tr}\left(Q_{\tau} D^{2}\right)$, whose spectrum is $(-\infty, 0]$. Take $\mu=a+i b \in \sigma\left(\mathcal{A}_{p}\right)$. The spectral inclusion Theorem IV.3.6 of [7] yields that $e^{\tau(a+i b)}$ belongs to $\sigma(T(\tau))=\sigma\left(T_{\tau}(1)\right)$. Since $\left(T_{\tau}(t)\right)_{t \geq 0}$ is analytic, from Corollary IV.3.12 of [7] we infer the identity $\sigma\left(T_{\tau}(1)\right) \backslash$ $\{0\}=e^{\sigma\left(A_{\tau}\right)}=(0,1]$. It follows that $a \leq 0$ and $\tau b=2 m \pi$ for some $m \in \mathbb{Z}$ and hence $i b \in \frac{2 \pi}{\tau} i \mathbb{Z}=G$, using also Proposition 2.3. Therefore $\sigma\left(\mathcal{A}_{p}\right)$ is contained in $(-\infty, 0]+G$.

The spectrum of the Ornstein-Uhlenbeck operators has been computed in [17, Section 4] if either $\sigma(B) \subset \mathbb{C}_{-}$or $\sigma(B) \subset \mathbb{C}_{+}$. The proofs in this paper are written only in the uniformly elliptic case where $Q$ is positive definite, but in the introduction of [17] it is pointed out that they also work only assuming the hypoellipticity condition (1.2).

To explain why this condition suffices, we recall that the spectrum of $\mathcal{A}_{p}$ is determined in [17] at first under the assumption $\sigma(B) \subset \mathbb{C}_{+}$by exhibiting explicit eigenfunctions for the eigenvalues $\mu<-\operatorname{tr}(B) / p$. These are computed using the matrix

$$
\widetilde{Q}_{\infty}=\int_{0}^{\infty} e^{-s B} Q e^{-s B^{T}} d s
$$

The above integral converges since the matrix semigroup $\left(e^{-s B}\right)_{s \geq 0}$ is exponentially stable. Moreover $\widetilde{Q}_{\infty}$ is nondegenerate under condition (1.2). Since $\widetilde{Q}_{\infty}$, and not $Q$, enters all calculations, all results still hold in the hypoelliptic setting provided that $\sigma(B) \subset \mathbb{C}_{+}$, including the extreme cases $p=1, \infty$.

The case $\sigma(B) \subset \mathbb{C}_{-}$follows from the preceding one by a simple duality argument, which we describe now. The formal adjoint of $\mathcal{A}$ is given by

$$
\mathcal{A}^{*}=\sum_{i, j=1}^{n} q_{i j} D_{i j}-\sum_{i, j=1}^{n} b_{i j} x_{j} D_{i}-\operatorname{tr}(B)
$$

Let $\mathcal{A}_{p^{\prime}}^{*}$ be the realization of $\mathcal{A}^{*}$ in $L^{p^{\prime}}, 1 / p+1 / p^{\prime}=1$, as the generator of the semigroup (1.3) with $-B$ replacing $B$ (also in the definition of $Q_{t}$ ) multiplied by the exponential factor $e^{-t \operatorname{tr}(B)}$. Notice that the spectrum of the drift matrix is now contained in $\mathbb{C}_{+}$. Therefore, for every $\mu \in \mathbb{C}$ with $\operatorname{Re} \mu<\operatorname{tr}(B) / p^{\prime}-\operatorname{tr}(B)=-\operatorname{tr}(B) / p$ the operator $\mu-\mathcal{A}_{p^{\prime}}^{*}$ is not injective. Let $\left(\mathcal{A}_{p}{ }^{\prime}, D\left(\mathcal{A}_{p}{ }^{\prime}\right)\right)$ denote the adjoint of $\mathcal{A}_{p}$ in $L^{p^{\prime}}$. Recalling that $C_{c}^{\infty}$ is a core for
$\mathcal{A}_{p^{\prime}}^{*}$ and $\mathcal{A}_{p}$, it is easily seen that $D\left(\mathcal{A}_{p^{\prime}}^{*}\right) \subseteq D\left(\mathcal{A}_{p}{ }^{\prime}\right)$ and $\mathcal{A}_{p}{ }^{\prime} f=\mathcal{A}_{p^{\prime}}^{*} f$ for every $f \in D\left(\mathcal{A}_{p^{\prime}}^{*}\right)$. Since $\mu-\mathcal{A}_{p^{\prime}}^{*}$ is not injective, it follows that $\mu-\mathcal{A}_{p}$ is not surjective and hence $\mu \in \sigma\left(\mathcal{A}_{p}\right)$. The other inclusion is provided by Corollary 2.7. Note that this works in the extreme cases $p=1, \infty$ as well.

We state the results discussed above.
Theorem 2.9. Let $1 \leq p \leq \infty$ and (1.2) be true. If either $\sigma(B) \subset \mathbb{C}_{-}$or $\sigma(B) \subset \mathbb{C}_{+}$, then $\sigma\left(\mathcal{A}_{p}\right)=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$. In the latter case, every $\mu$ with $\operatorname{Re} \mu<-\operatorname{tr}(B) / p$ is an eigenvalue.

## 3 Preliminary considerations

In contrast to [17] we cannot use eigenvalues in the proof of our main result. To show this we rule out eigenvalues of $\mathcal{A}$ if the spectrum of $B$ intersects $\mathbb{C}_{-}$or if $B$ is skew-symmetric and $Q=I$, where we assume that $p=2$.

First, we assume that some eigenvalue of $B$ has a negative real part. Suppose that $\mu \in \mathbb{C}$ with $\operatorname{Re} \mu<-\frac{1}{2} \operatorname{tr}(B)$ was an eigenvalue of $\mathcal{A}_{2}$ with eigenfunction $f \in L^{2} \backslash\{0\}$. The spectral mapping theorem for the point spectrum shows that $T(t) f=e^{\mu t} f$ for every $t \geq 0$, see Theorem IV.3.7 and Corollary IV.3.8 of [7]. Denoting by $\hat{f}$ the Fourier transform of $f$, the representation formula (1.3) implies that the equation $T(t) f=e^{\mu t} f$ is equivalent to

$$
\begin{equation*}
\hat{f}\left(e^{-t B^{T}} \xi\right)=e^{(\mu+\operatorname{tr}(B)) t} e^{\left|Q_{t}^{1 / 2} e^{-t B^{T}} \xi\right|^{2}} \hat{f}(\xi) \tag{3.1}
\end{equation*}
$$

where $t \geq 0$, see Section 4 of [17]. We compute

$$
\begin{align*}
\left|Q_{t}^{1 / 2} e^{-t B^{T}} \xi\right|^{2} & =\left\langle Q_{t} e^{-t B^{T}} \xi, e^{-t B^{T}} \xi\right\rangle=\int_{0}^{t}\left\langle e^{s B} Q e^{s B^{T}} e^{-t B^{T}} \xi, e^{-t B^{T}} \xi\right\rangle d s \\
& =\int_{0}^{t}\left|Q^{1 / 2} e^{(s-t) B^{T}} \xi\right|^{2} d s=\int_{0}^{t}\left|Q^{1 / 2} e^{-s B^{T}} \xi\right|^{2} d s \tag{3.2}
\end{align*}
$$

for $\xi \in \mathbb{R}^{n}$. Take $\lambda \in \sigma(B)=\sigma\left(B^{T}\right)$ with $\operatorname{Re} \lambda<0$. Let $P$ be the spectral projection of $B^{T}$ corresponding to $\lambda$. Fix $\varepsilon>0$ with $\operatorname{Re} \lambda+\varepsilon<0$. Then there exists a constant $M>0$ such that $\left\|e^{s B^{T}} P\right\| \leq M e^{(\operatorname{Re} \lambda+\varepsilon) s}$ for every $s \geq 0$. Observe that also $-B$ satisfies (1.2), so that there is a constant $\nu>0$ with

$$
\int_{0}^{1}\left|Q^{1 / 2} e^{-s B^{T}} \xi\right|^{2} d s=\left\langle\int_{0}^{1} e^{-s B} Q e^{-s B^{T}} \xi d s, \xi\right\rangle \geq \nu|\xi|^{2}
$$

Let $t \in[m, m+1)$ for some $m \in \mathbb{N}_{0}$. Inserting $P$ in (3.2), it follows

$$
\begin{aligned}
\left|Q_{t}^{1 / 2} e^{-t B^{T}} \xi\right|^{2} & \geq \int_{0}^{m}\left|Q^{1 / 2} e^{-s B^{T}} \xi\right|^{2} d s=\sum_{k=0}^{m-1} \int_{0}^{1}\left|Q^{1 / 2} e^{-r B^{T}} e^{-k B^{T}} \xi\right|^{2} d r \\
& \geq \sum_{k=0}^{m-1} \nu\|P\|^{-2}\left|P e^{-k B^{T}} \xi\right|^{2} \geq \sum_{k=0}^{m-1} \nu(M\|P\|)^{-2} e^{-2(\operatorname{Re} \lambda+\varepsilon) k}|P \xi|^{2} \\
& \geq c e^{-2(\operatorname{Re} \lambda+\varepsilon) t}|P \xi|^{2}
\end{aligned}
$$

for some constant $c>0$. Integrating (3.1) on $\mathbb{R}^{n}$, we derive

$$
\begin{aligned}
\int_{\mathbb{R}^{n}}|\hat{f}(\xi)|^{2} d \xi & =e^{2 t\left(\operatorname{Re} \mu+\frac{1}{2} \operatorname{tr}(B)\right)} \int_{\mathbb{R}^{n}} e^{2\left|Q_{t}^{1 / 2} e^{-t B^{T}} \xi\right|^{2}}|\hat{f}(\xi)|^{2} d \xi \\
& \geq \exp \left(2 c \alpha^{2} e^{-2(\operatorname{Re} \lambda+\varepsilon) t}\right) e^{2 t\left(\operatorname{Re} \mu+\frac{1}{2} \operatorname{tr}(B)\right)} \int_{\{|P \xi| \geq \alpha\}}|\hat{f}(\xi)|^{2} d \xi
\end{aligned}
$$

for every $t \geq 0$ and $\alpha>0$. Letting $t \rightarrow+\infty$, the right hand side blows up unless $\int_{\{|P \xi| \geq \alpha\}}|\hat{f}(\bar{\xi})|^{2} d \xi=0$. Since $\alpha>0$ is arbitrary and the set $\{P \xi=0\}$ has measure 0 , this would imply $\hat{f}=0$ and thus $f=0$ in $L^{2}$, which is a contradiction.

Second, we assume that $B=-B^{T}$ and $Q=I$. Recalling that $\operatorname{tr}(B)=0$, we now suppose there was an eigenvalue $\mu$ of $\mathcal{A}_{2}$ with $\operatorname{Re} \mu<0$. Arguing as before, we rewrite (3.1) as

$$
e^{-\left|Q_{t}^{1 / 2} e^{-t B^{T}} \xi\right|^{2}} \hat{f}\left(e^{-t B^{T}} \xi\right)=e^{\mu t} \hat{f}(\xi), \quad t \geq 0
$$

and then integrate over $\mathbb{R}^{n}$ to obtain

$$
\int_{\mathbb{R}^{n}} e^{-2\left|Q_{t}^{1 / 2} \xi\right|^{2}}|\hat{f}(\xi)|^{2} d \xi=e^{2 \operatorname{Re} \mu t} \int_{\mathbb{R}^{n}}|\hat{f}(\xi)|^{2} d \xi
$$

Observing

$$
\left|Q_{t}^{1 / 2} \xi\right|^{2}=\left\langle Q_{t} \xi, \xi\right\rangle=\int_{0}^{t}\left\langle e^{s B} e^{s B^{T}} \xi, \xi\right\rangle d s=\int_{0}^{t}\left\langle e^{s\left(B+B^{T}\right)} \xi, \xi\right\rangle d s=t|\xi|^{2}
$$

we derive

$$
\begin{aligned}
\int_{\mathbb{R}^{n}}|\hat{f}(\xi)|^{2} d \xi & =\int_{\mathbb{R}^{n}} e^{-2 t\left(|\xi|^{2}+\operatorname{Re} \mu\right)}|\hat{f}(\xi)|^{2} d \xi \\
& =\int_{\left\{|\xi|^{2}>-\operatorname{Re} \mu\right\}} e^{-2 t\left(|\xi|^{2}+\operatorname{Re} \mu\right)}|\hat{f}(\xi)|^{2} d \xi+\int_{\left\{|\xi|^{2}<-\operatorname{Re} \mu\right\}} e^{-2 t\left(|\xi|^{2}+\operatorname{Re} \mu\right)}|\hat{f}(\xi)|^{2} d \xi
\end{aligned}
$$

The first integral in the last line tends to 0 as $t \rightarrow+\infty$ by dominated convergence. The second integral tends to $+\infty$ by monotone convergence, if $\int_{\left\{|\xi|^{2}<-\operatorname{Re} \mu\right\}}|\hat{f}(\xi)|^{2} d \xi>0$. Therefore we have either $\|\hat{f}\|_{2}=+\infty$ or $\|\hat{f}\|_{2}=0$, and we get a contradiction in any case.

By duality one deduces from the above examples that, if $\sigma(B)$ intersects both $\mathbb{C}_{-}$and $\mathbb{C}_{+}$, a point $\lambda$ can be in the spectrum of $\mathcal{A}_{2}$ even though $\lambda-\mathcal{A}_{2}$ is injective and has dense range. Approximate eigenvalues will thus play a central role.

In order to describe the spectrum of $\mathcal{A}$, we will reduce the drift matrix $B$ to a quasi triangular upper matrix. This is done as follows. If $M$ is an invertible real $n \times n$ matrix, we define the change of variables

$$
\begin{equation*}
\Phi_{M}: L^{p} \rightarrow L^{p}, \quad\left(\Phi_{M} u\right)(y)=u\left(M^{-1} y\right) \tag{3.3}
\end{equation*}
$$

Setting $v=\Phi_{M} u$, one easily calculates that $\mathcal{A} u(x)=\mathcal{A}_{0} v(M x)$ for $x \in \mathbb{R}^{n}$, where

$$
\mathcal{A}_{0} v=\operatorname{Tr}\left(Q_{0} D^{2} v\right)+\left\langle B_{0} y, D v\right\rangle
$$

with $y=M x, Q_{0}=M Q M^{T}$, and $B_{0}=M B M^{-1}$. We conclude

$$
\mathcal{A}=\Phi_{M}^{-1} \mathcal{A}_{0} \Phi_{M} \quad \text { with } \quad D\left(\mathcal{A}_{p}\right)=\Phi_{M}^{-1} D\left(\mathcal{A}_{0, p}\right)
$$

We observe that the new operator $\mathcal{A}_{0}$ is still hypoelliptic, see (1.2), and that the spectrum is invariant under this transformation.

Applying Schur's theorem for real matrices (see e.g. Theorem 2.3.4 in [11]), we can now choose a real orthogonal matrix $M$ such that $M B M^{-1}=T$ with

$$
T=\left(\begin{array}{ccccc}
B_{1} & * & * & \cdots & *  \tag{3.4}\\
0 & B_{2} & * & \cdots & * \\
0 & 0 & B_{3} & * & \vdots \\
\vdots & \vdots & & \ddots & \vdots \\
0 & 0 & 0 & \cdots & B_{l}
\end{array}\right)
$$

where each $B_{j}$ is a real $1 \times 1$ matrix with a real eigenvalue of $B$, or a real $2 \times 2$ matrix with a pair of nonreal complex conjugate eigenvalues $\alpha_{j} \pm i \beta_{j}$. The diagonal blocks $B_{j}$ may be arranged in any prescribed order. By * we denote an arbitrary block.

## 4 The spectrum of $\mathcal{A}_{p}$

The spectrum of the Ornstein-Uhlenbeck operators $\mathcal{A}_{p}$ depends on the spectrum of the drift operator $\mathcal{L}_{p}$ which in turn is determined by $B$. If $B$ is similar to a diagonal matrix with nonzero eigenvalues $\pm i \sigma_{1}, \pm i \sigma_{2}, \ldots, \pm i \sigma_{k}$ in $i \mathbb{R}$ and possibly 0 and if all ratios $\sigma_{r} \sigma_{s}^{-1}$ belong to $\mathbb{Q}$, then Theroem 2.2 shows that $\sigma\left(\mathcal{L}_{p}\right)$ is a discrete subgroup $G=\frac{2 \pi i}{\tau} \mathbb{Z}$ of $i \mathbb{R}$, where $S(\tau)=I$. In this case we prove that $\sigma\left(\mathcal{A}_{p}\right)=(-\infty, 0]+G$. In all the remaining cases, the spectrum of $\mathcal{L}_{p}$ is the vertical line $-\operatorname{tr}(B) / p+i \mathbb{R}$ and we show that $\sigma\left(\mathcal{A}_{p}\right)$ is the half-plane $\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$. These results, which are the main achievement of the paper, are stated in Theorem 1.1, which we rewrite below for convenience.

Theorem 4.1. Let (1.2) be true and $p \in[1, \infty]$. Then the spectrum of $\mathcal{A}_{p}$ is given by

$$
\begin{equation*}
\sigma\left(\mathcal{A}_{p}\right)=(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \tag{4.1}
\end{equation*}
$$

In particular, we have either $\sigma\left(\mathcal{A}_{p}\right)=(-\infty, 0]+G$ or $\sigma\left(\mathcal{A}_{p}\right)=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$, according to $\sigma\left(\mathcal{L}_{p}\right)$ being the discrete subgroup $G=\frac{2 \pi i}{\tau} \mathbb{Z}$ of $i \mathbb{R}$ or the whole line $-\operatorname{tr}(B) / p+$ $i \mathbb{R}$. Moreover, the semigroup $(T(t))_{t \geq 0}$ satisfies the weak spectral mapping theorem

$$
\begin{equation*}
\sigma(T(t))=\overline{\exp \left(t \sigma\left(\mathcal{A}_{p}\right)\right)}, \quad t \geq 0 \tag{4.2}
\end{equation*}
$$

We even have $\sigma(T(t)) \backslash\{0\}=e^{t \sigma\left(\mathcal{A}_{p}\right)}$ except for the case that $\sigma\left(\mathcal{L}_{p}\right)=G=\frac{2 \pi i}{\tau} \mathbb{Z}$ and $t / \tau$ is irrational.

Proof. Theorem 2.9 shows the equality (4.1) if $\sigma(B) \subset \mathbb{C}_{-}$or $\sigma(B) \subset \mathbb{C}_{+}$. Moreover, by Corollary 2.7 and Proposition 2.8 the inclusion $\sigma\left(\mathcal{A}_{p}\right) \subseteq(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$ always holds. Therefore we only have to prove the other inclusion $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$ in two remaining cases: one eigenvalue of $B$ has a positive real part or one eigenvalue of $B$ lies on the imaginary axis. Note that these cases may overlap and that the first one includes situations covered by Theorem 2.9. The inclusion $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$ is established in these two cases in the following two subsections.

To prove the (weak) spectral mapping theorem, we take (4.1) for granted. Let $t>0$. The spectral inclusion Theorem IV.3.6 of [7] and Corollary 2.7 show that

$$
e^{t \sigma\left(\mathcal{A}_{p}\right)} \subseteq \sigma(T(t)) \backslash\{0\} \subseteq\left\{\mu \in \mathbb{C}\left|0<|\mu| \leq e^{-t \operatorname{tr}(B) / p}\right\}=: B_{t}\right.
$$

We thus even obtain $\sigma(T(t)) \backslash\{0\}=e^{\operatorname{t\sigma }\left(\mathcal{A}_{p}\right)}$ if $\sigma\left(\mathcal{A}_{p}\right)=\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu \leq-\operatorname{tr}(B) / p\}$. In the other case Proposition 2.3 yields $\sigma\left(\mathcal{L}_{p}\right)=G=\frac{2 \pi i}{\tau} \mathbb{Z}$. Now $\operatorname{tr}(B)=0$. By (4.1) we can now write

$$
e^{t \sigma\left(\mathcal{A}_{p}\right)}=\left\{e^{t a} e^{2 \pi i m t / \tau} \mid a \leq 0, m \in \mathbb{Z}\right\}=(0,1] \cdot S_{t}
$$

with $S_{t}=\left\{e^{2 \pi i m t / \tau} \mid m \in \mathbb{Z}\right\}$. There are two subcases.
First, let $t / \tau$ be irrational. Then the set $S_{t}$ is dense in the unit circle and it follows that $\overline{\exp \left(t \sigma\left(\mathcal{A}_{p}\right)\right)}$ is equal to $\overline{B_{0}}$; i.e., (4.2) is true.

Second, let $t / \tau=j / k$ for coprime $j, k \in \mathbb{N}$. Then $S_{t}$ coincides with the set $\Gamma_{k}$ of $k$ th unit roots so that $e^{t \sigma\left(\mathcal{A}_{p}\right)}=(0,1] \cdot \Gamma_{k}$ On the other hand, we have $S(t)^{k}=S(j \tau)=I$. As in the proof of Proposition 2.8, we deduce that $T(t)^{k}=T(k t)=T_{k t}(1)$ for the analytic semigroup $\left(T_{k t}(s)\right)_{s \geq 0}$ generated by $\operatorname{Tr}\left(Q_{k t} D^{2}\right)$. The spectrum of $T_{k t}(1)$ is thus equal to $[0,1]$ and hence $\sigma(T(t)) \backslash\{0\}=(0,1] \cdot \Gamma_{k}$ as required.

### 4.1 The case $\sigma(B) \cap \mathbb{C}_{+} \neq \emptyset$

We show the remaining inclusion in the proof of Theorem 4.1 in the first case.
Proposition 4.2. Let $\sigma(B) \cap \mathbb{C}_{+} \neq \emptyset$. Then $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$.
In the proof we use degenerate Ornstein-Uhlenbeck operators depending on different sets of variables, as we explain now. We let $\mathbb{R}^{n}=\mathbb{R}^{k} \times \mathbb{R}^{m}$ and write a point $z \in \mathbb{R}^{n}$ accordingly as $z=(x, y)$. Let $B_{1}$ and $B_{2}$ be real $k \times k$ and $m \times m$ matrices, respectively, and $Q_{2}$ a real, symmetric and positive semidefinite $m \times m$ matrix. We consider the operators

$$
\begin{equation*}
\mathcal{L}^{(1)}=\left\langle B_{1} x, D_{x}\right\rangle \quad \text { and } \quad \mathcal{A}^{(2)}=\operatorname{Tr}\left(Q_{2} D^{2}\right)+\left\langle B_{2} y, D_{y}\right\rangle \tag{4.3}
\end{equation*}
$$

Here $\mathcal{L}^{(1)}$ is a drift operator on $L^{p}\left(\mathbb{R}^{k}\right)$ and $\mathcal{A}^{(2)}$ is an Ornstein-Uhlenbeck operator on $L^{p}\left(\mathbb{R}^{m}\right)$, which is assumed to be hypoelliptic (recall that $L^{\infty}$ means $C_{0}$ ). Let $\left(S_{1}(t)\right)_{t \geq 0}$ and $\left(T_{2}(t)\right)_{t \geq 0}$ be the generated semigroups. Then $\left(S_{1}(t) \otimes T_{2}(t)\right)_{t \geq 0}$ acting on $L^{p}\left(\mathbb{R}^{k}\right) \otimes L^{p}\left(\mathbb{R}^{m}\right)$ can be extended to a $C_{0}$-semigroup on $L^{p}\left(\mathbb{R}^{n}\right)$, whose generator is the closure $\mathcal{C}_{p}$ of $\mathcal{C}=$ $\mathcal{L}^{(1)}+\mathcal{A}^{(2)}$ initially defined on $D\left(\mathcal{L}_{p}^{(1)}\right) \otimes D\left(\mathcal{A}_{p}^{(2)}\right)$. Since $C_{c}^{\infty}\left(\mathbb{R}^{k}\right)$ and $C_{c}^{\infty}\left(\mathbb{R}^{m}\right)$ are cores for $\mathcal{L}^{(1)}$ and $\mathcal{A}^{(2)}$, respectively, it follows that $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ is a core for $\mathcal{C}_{p}$.

Proof of Proposition 4.2. Let $\lambda=\alpha+i \beta \in \sigma(B)$ with $\alpha>0$. As explained at the end of Section 3, using a change of variables we can assume that our operator is given by

$$
\mathcal{A}=\operatorname{Tr}\left(Q_{0} D^{2}\right)+\langle T x, D\rangle
$$

where $T$ is in the quasi triangular form (3.4), its last block $B_{l}$ corresponds to $\lambda$, and $Q_{0}$ is the transformed diffusion matrix. We distinguish between the cases $\beta=0$ and $\beta \neq 0$. (Below we tacitly assume that $T \neq B_{l}$ since the easier case $T=B_{l}$ can be treated analogously.)

Case 1. $\boldsymbol{\beta}=\mathbf{0}$. Denote the nonreal eigenvalues of $B$ by $\left\{\alpha_{1} \pm i \beta_{1}, \ldots, \alpha_{k} \pm i \beta_{k}\right\}$ with $0 \leq 2 k<n$ and the real ones by $\left\{\lambda_{2 k+1}, \ldots, \lambda_{n}\right\}$ with $\lambda_{n}=\alpha>0$. We write a point in $\mathbb{R}^{n}$ as $x=\left(x_{1}, y_{1}, \ldots, x_{k}, y_{k}, w_{2 k+1}, \ldots, w_{n}\right)$. We use a scaling argument in which the variables $z_{j}=\left(x_{j}, y_{j}\right)$ relative to conjugated eigenvalues are coupled and which leaves the last variable unscaled. Let $D_{z_{j}}=\left(D_{x_{j}}, D_{y_{j}}\right)$ for $j=1, \ldots, k$. The scaling operator is defined by

$$
I_{r} u(x)=u\left(\frac{z_{1}}{r^{\gamma_{1}}}, \frac{z_{2}}{r^{\gamma_{2}}}, \ldots, \frac{z_{k}}{r^{\gamma_{k}}}, \frac{w_{2 k+1}}{r^{\gamma_{2 k+1}}}, \ldots, \frac{w_{n-1}}{r^{\gamma_{n-1}}}, w_{n}\right)
$$

for $r>0$ and with $\gamma_{1}=1$ and $\gamma_{i}>\gamma_{j}>0$ for $i<j$. Observe that $\left\|I_{r}^{-1}\right\|=\left\|I_{r}\right\|^{-1}$ on $L^{p}$. Let $u \in C_{c}^{\infty}$. Computing $I_{r}^{-1} \mathcal{A} I_{r} u$, one finds that

$$
\lim _{r \rightarrow+\infty} I_{r}^{-1} \mathcal{A} I_{r} u=\mathcal{C} u \quad \text { in } L^{p},
$$

for the limit operator

$$
\mathcal{C} u=\nu D_{w_{n}}^{2} u+\lambda_{n} w_{n} D_{w_{n}} u+\sum_{j=1}^{k}\left\langle B_{j} z_{j}, D_{z_{j}} u\right\rangle+\sum_{j=2 k+1}^{n-1} \lambda_{j} w_{j} D_{w_{j}} u .
$$

The constant $\nu$ is the component $\left\langle Q_{0} e_{n}, e_{n}\right\rangle$ of $Q_{0}$ where $e_{n}=(0, \ldots, 0,1)$. It is positive, which can be explained as follows. The last row vector in the matrix $T$ is $\lambda_{n} e_{n}$. This means that the transpose of $T$ maps $e_{n}$ to $\lambda_{n} e_{n}$. Let $X$ be the one-dimensional subspace spanned by $e_{n}$. It is invariant for the transpose of $T$. Since $\mathcal{A}$ is hypoelliptic, $X$ is not contained in the kernel of $Q_{0}$. It follows $Q_{0} e_{n} \neq 0$ and hence $\nu=\left|Q_{0}^{1 / 2} e_{n}\right|^{2}>0$.

Note that we can write $\mathcal{C}=\mathcal{L}^{\text {lim }}+\mathcal{A}^{\text {lim }}$ with

$$
\mathcal{A}^{\lim }=\nu D_{n}^{2}+\lambda_{n} w_{n} D_{n}, \quad \mathcal{L}^{\lim }=\sum_{j=1}^{k}\left\langle B_{j} z_{j}, D_{z_{j}}\right\rangle+\sum_{j=2 k+1}^{n-1} \lambda_{j} w_{j} D_{w_{j}} .
$$

We endow $\mathcal{C}$ with the domain described before the proof and call it $\mathcal{C}_{p}$. We first establish a crucial spectral property of $\mathcal{C}_{p}$.

Claim. Every $\mu$ with $\operatorname{Re} \mu<-\operatorname{tr}(B) / p$ is an approximate eigenvalue for $\mathcal{C}_{p}$.
Since $\lambda_{n}>0$, every $\lambda$ with $\operatorname{Re} \lambda<-\lambda_{n} / p$ is an eigenvalue of the one dimensional operator $\mathcal{A}_{p}^{\lim }$ by Theorem 2.9. Theorem 2.2 shows that $\mathcal{L}_{p}^{\lim }$ possesses the approximate eigenvalue $-c / p$, where

$$
c=2 \sum_{j=1}^{k} \alpha_{j}+\sum_{i=2 k+1}^{n-1} \lambda_{i}=\operatorname{tr}(B)-\lambda_{n} .
$$

Now, fix $\mu \in \mathbb{C}$ with $\operatorname{Re} \mu<-\operatorname{tr}(B) / p$ and set $\lambda=\mu+c / p$. Note that $\operatorname{Re} \lambda<-\lambda_{n} / p$. Choose an eigenfunction $u_{1}=u_{1}\left(w_{n}\right)$ of $\mathcal{A}_{p}^{\lim }$ for $\lambda$ with $\left\|u_{1}\right\|_{L^{p}(\mathbb{R})}=1$. Given $\varepsilon>0$, there is a function $u_{2}=u_{2}\left(z_{1}, \ldots, z_{k}, w_{2 k+1}, \ldots, w_{n-1}\right)$ in $D\left(\mathcal{L}_{p}^{\lim }\right)$ with $\left\|u_{2}\right\|_{L^{p}\left(\mathbb{R}^{n-1}\right)}=1$ and $\left\|\mathcal{L}^{\lim } u_{2}+\frac{c}{p} u_{2}\right\|_{L^{p}\left(\mathbb{R}^{n-1}\right)} \leq \varepsilon$. The function $u=u_{1} u_{2}$ thus belongs to $D\left(\mathcal{C}_{p}\right)$, has norm one in $L^{p}$ and satisfies

$$
\begin{equation*}
\mathcal{C} u-\mu u=\left(\mathcal{A}^{\lim } u_{1}-\lambda u_{1}\right) u_{2}+\left(\mathcal{L}^{\lim } u_{2}+\frac{c}{p} u_{2}\right) u_{1}=\left(\mathcal{L}^{\lim } u_{2}+\frac{c}{p} u_{2}\right) u_{1}, \tag{4.4}
\end{equation*}
$$

which yields $\|\mathcal{C} u-\mu u\|_{p} \leq \varepsilon$. So the claim is proved.
Take $\lambda_{0} \in \rho\left(\mathcal{A}_{p}\right)$. By similarity, $\lambda_{0}$ belongs to $\rho\left(I_{r}^{-1} \mathcal{A} I_{r}\right)$ with resolvent $R\left(\lambda_{0}, I_{r}^{-1} \mathcal{A}_{p} I_{r}\right)=$ $I_{r}^{-1} R\left(\lambda_{0}, \mathcal{A}_{p}\right) I_{r}$ for all $r>0$. It follows $\left\|R\left(\lambda_{0}, I_{r}^{-1} \mathcal{A}_{p} I_{r}\right)\right\| \leq \| R\left(\lambda_{0}, \mathcal{A}_{p} \|=: C\right.$ and

$$
\begin{equation*}
\|u\|_{p}=\left\|R\left(\lambda_{0}, I_{r}^{-1} \mathcal{A}_{p} I_{r}\right)\left(\lambda_{0}-I_{r}^{-1} \mathcal{A} I_{r}\right) u\right\|_{p} \leq C\left\|\left(\lambda_{0}-I_{r}^{-1} \mathcal{A} I_{r}\right) u\right\|_{p} \tag{4.5}
\end{equation*}
$$

for $u \in C_{c}^{\infty}$. Letting $r \rightarrow+\infty$, we infer $\|u\|_{p} \leq C\left\|\left(\lambda_{0}-\mathcal{C}\right) u\right\|_{p}$. Since $C_{c}^{\infty}$ is a core for $\mathcal{C}_{p}$, this shows that $\lambda_{0}$ cannot be an approximate eigenvalue of $\mathcal{C}_{p}$, and hence $\operatorname{Re} \lambda_{0} \geq-\operatorname{tr}(B) / p$ by the claim. This means that $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$ is contained in $\sigma\left(\mathcal{A}_{p}\right)$.

Case 2. $\boldsymbol{\beta} \neq \mathbf{0}$. We rearrange the blocks in (3.4) such that the first blocks contain the real eigenvalues $\lambda_{1}, \ldots, \lambda_{k}$ with $0 \leq k<n$ and the other blocks contain the complex ones
$\alpha_{k+1} \pm i \beta_{k+1}, \ldots, \alpha_{k+m} \pm i \beta_{k+m}$ for $n=k+2 m$, where $\alpha_{k+m}=\alpha>0$ and $\beta_{k+m}=\beta \neq 0$. As a consequence, a point of $\mathbb{R}^{n}$ is denoted by $x=\left(w_{1}, w_{2}, \ldots, w_{k}, z_{k+1}, \ldots, z_{k+m}\right)$ with $z_{j}=\left(x_{j}, y_{j}\right)$. The scaling operator is now defined by

$$
\begin{equation*}
J_{r} u(x)=u\left(\frac{w_{1}}{r^{\gamma_{1}}}, \frac{w_{2}}{r^{\gamma_{2}}}, \ldots, \frac{w_{k}}{r^{\gamma_{k}}}, \frac{z_{k+1}}{r^{\gamma_{k+1}}}, \ldots, \frac{z_{k+m-1}}{r^{\gamma_{k+m-1}}}, z_{k+m}\right) \tag{4.6}
\end{equation*}
$$

with $\gamma_{1}=1$ and $\gamma_{i}>\gamma_{j}>0$ for $i<j$. For every $u \in C_{c}^{\infty}$ we have

$$
\lim _{r \rightarrow+\infty} J_{r}^{-1} \mathcal{A} J_{r} u=\mathcal{C} u \quad \text { in } L^{p}
$$

where the limit operator is given by

$$
\begin{align*}
\mathcal{C} u=\operatorname{Tr} & \left(Q_{0}^{\dagger} D_{k+m}^{2} u\right)+\left\langle B_{k+m} z_{k+m}, D_{z_{k+m}} u\right\rangle  \tag{4.7}\\
& +\sum_{j=1}^{k} \lambda_{j} w_{j} D_{w_{j}} u+\sum_{j=1}^{m-1}\left\langle B_{k+j} z_{k+j}, D_{z_{k+j}} u\right\rangle .
\end{align*}
$$

Here $Q_{0}^{\dagger}$ is the lower right $2 \times 2$ submatrix of $Q_{0}$.
As before we introduce $\mathcal{C}_{p}$ and claim that the open half-plane $\{\mu \in \mathbb{C} \mid \operatorname{Re} \mu<-\operatorname{tr}(B) / p\}$ is contained in its approximate spectrum.

To prove the claim, we split $\mathcal{C}$ as the sum $\mathcal{A}^{\lim }+\mathcal{L}^{\text {lim }}$ for the Ornstein-Uhlenbeck operator

$$
\mathcal{A}^{\lim }=\operatorname{Tr}\left(Q_{0}^{\dagger} D_{k+m}^{2}\right)+\left\langle B_{k+m} z_{k+m}, D_{z_{k+m}}\right\rangle
$$

in the last two variables and the drift operator

$$
\mathcal{L}^{\lim }=\sum_{j=1}^{k} \lambda_{j} w_{j} D_{w_{j}}+\sum_{j=1}^{m-1}\left\langle B_{k+j} z_{k+j}, D_{z_{k+j}}\right\rangle
$$

acting in the remaining variables.
We show that $\mathcal{A}^{\text {lim }}$ is hypoelliptic by verifying (1.2). Let $Y$ be a real subspace of $\mathbb{R}^{2}$ which is invariant for $B_{k+m}^{T}$. Suppose that $\operatorname{dim} Y=1$. Then there would exist a real eigenvalue for $B_{k+m}^{T}$, but this is not the case as $\sigma\left(B_{k+m}^{T}\right)=\{\alpha \pm i \beta\}$. We thus have either $Y=\{0\}$ or $Y=\mathbb{R}^{2}$. Suppose that $\mathbb{R}^{2} \subseteq \operatorname{ker}\left(Q_{0}^{\dagger}\right)$. In this case the real subspace of $\mathbb{R}^{n}$ spanned by $e_{n-1}=(0, \ldots, 0,1,0)$ and $e_{n}=(0, \ldots, 0,1)$ would be invariant for the transpose of the drift matrix $T$ and it would be contained in $\operatorname{ker}\left(Q_{0}\right)$. This contradicts the hypoellipticity of $\mathcal{A}$.

Since $\alpha_{k+m}>0$, by Theorem 2.9 every $\lambda$ with $\operatorname{Re} \lambda<-2 \alpha_{k+m} / p$ is an eigenvalue of $\mathcal{A}_{p}^{\lim }$. Moreover, Theorem 2.2 yields that $\sigma_{a p}\left(\mathcal{L}_{p}^{\lim }\right)$ contains the number $-c / p$ with

$$
c=\sum_{j=1}^{k} \lambda_{j}+\sum_{j=1}^{m-1} 2 \alpha_{j+k} .
$$

Take $\mu \in \mathbb{C}$ with $\operatorname{Re} \mu<-\operatorname{tr}(B) / p$ and set $\lambda=\mu+c / p$. Take an eigenfunction $u_{1}=$ $u_{1}\left(x_{k+m}, y_{k+m}\right)$ of $\mathcal{A}_{p}^{\text {lim }}$ for $\lambda$ with $\left\|u_{1}\right\|_{L^{p}\left(\mathbb{R}^{2}\right)}=1$. Given $\varepsilon>0$, we have a function $u_{2}=u_{2}\left(w_{1}, \ldots, w_{k}, z_{k+1}, \ldots, z_{k+m-1}\right)$ in $D\left(\mathcal{L}_{p}^{\lim }\right)$ satisfying $\left\|u_{2}\right\|_{L^{p}\left(\mathbb{R}^{n-2}\right)}=1$ and $\| \mathcal{L}^{\lim } u_{2}+$ $\frac{c}{p} u_{2} \|_{L^{p}\left(\mathbb{R}^{n-2}\right)} \leq \varepsilon$. As in (4.4), $u=u_{1} u_{2}$ is an approximate eigenfunction for $\mathcal{C}_{p}$ with approximate eigenvalue $\mu$. We can then proceed as at the end of Case 1, see (4.5).

### 4.2 The case $\sigma(B) \cap i \mathbb{R} \neq \emptyset$

To deal with imaginary eigenvalues of $B$, we need a second type of transformation. We introduce an isometry $S: L^{p} \rightarrow L^{p}$ by

$$
\begin{equation*}
S u(x)=e^{i s \phi(x)} u(x) \tag{4.8}
\end{equation*}
$$

where $s \in \mathbb{R}$ and the function $\phi: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is chosen below. For $u \in C_{c}^{\infty}$, say, the operator $\mathcal{A}$ given by (1.1) is transformed into

$$
\begin{equation*}
S^{-1} \mathcal{A} S u(x)=\mathcal{A} u-s^{2}\langle Q D \phi, D \phi\rangle u+2 i s\langle Q D \phi, D u\rangle+i s \operatorname{Tr}\left(Q D^{2} \phi\right) u+i s\langle B x, D \phi\rangle u \tag{4.9}
\end{equation*}
$$

Let $\lambda \in \sigma(B) \cap i \mathbb{R}$. Then we have either $\lambda=0$ or $\lambda=i \beta \neq 0$. The next two propositions show the spectral inclusion needed for Theorem 4.1 separately for these two cases.
Proposition 4.3. Let $0 \in \sigma(B)$. Then $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$.
Proof. In the proof we write $\mathcal{A}$ for $\mathcal{A}_{p}$ and similarly for the other operators involved. Observe that the kernel of $B^{T}$ is a nontrivial subspace which is invariant for $B^{T}$. Condition (1.2) thus yields a vector $\xi \in \operatorname{ker}\left(B^{T}\right)$ with $Q \xi \neq 0$. Then $\langle Q \xi, \xi\rangle=\left|Q^{1 / 2} \xi\right|^{2} \neq 0$. We set $\phi(x)=\xi \cdot x$ for $x \in \mathbb{R}^{n}$. Equation (4.9) then becomes

$$
\begin{equation*}
S^{-1} \mathcal{A} S u(x)=\mathcal{A} u-s^{2}\langle Q \xi, \xi\rangle u+2 i s\langle Q \xi, D u\rangle+i s\left\langle x, B^{T} \xi\right\rangle u=\tilde{\mathcal{A}} u-s^{2}\langle Q \xi, \xi\rangle u \tag{4.10}
\end{equation*}
$$

where we have defined $\tilde{\mathcal{A}} u=\mathcal{A} u+2 i s\langle Q \xi, D u\rangle$ and used $B^{T} \xi=0$. Let $k \in \mathbb{N}$ and the isometry $V_{k}: L^{p} \rightarrow L^{p}$ be given by

$$
\begin{equation*}
V_{k} u(x)=k^{-n / p} u\left(k^{-1} x\right) \tag{4.11}
\end{equation*}
$$

For $u \in C_{c}^{\infty}$, we compute

$$
V_{k}^{-1} \tilde{\mathcal{A}} V_{k} u=k^{-2} \operatorname{Tr}\left(Q D^{2} u\right)+\langle B x, D u\rangle+k^{-1} 2 i s\langle Q \xi, D u\rangle \longrightarrow \mathcal{L} u=\langle B x, D u\rangle
$$

as $k \rightarrow+\infty$. Set $\tilde{\mathcal{A}}_{k}=V_{k}^{-1} \tilde{\mathcal{A}} V_{k}$. Then $\rho\left(\tilde{\mathcal{A}}_{k}\right)=\rho(\tilde{\mathcal{A}})$ by similarity, where we omit the subscript $p$. We want to show the inclusion $\sigma(\mathcal{L}) \subseteq \sigma(\tilde{\mathcal{A}})$, for which we need the next fact.

Claim. Let $\lambda \in \rho(\tilde{\mathcal{A}}) \cap \rho(\mathcal{L})$ and $f \in L^{p}$. We then obtain

$$
\begin{equation*}
R\left(\lambda, \tilde{\mathcal{A}}_{k}\right) f \rightarrow R(\lambda, \mathcal{L}) f \quad \text { in } L^{p} \quad \text { as } k \rightarrow \infty \tag{4.12}
\end{equation*}
$$

Since $C_{c}^{\infty}$ is a core of $\left(\mathcal{L}, D\left(\mathcal{L}_{p}\right)\right)$ by Proposition 2.1, it suffices to prove the convergence on the dense subspace $(\lambda-\mathcal{L}) C_{c}^{\infty}$. Let $f=\lambda u-\mathcal{L} u$ for some $u \in C_{c}^{\infty}$. Using the identity

$$
R\left(\lambda, \tilde{\mathcal{A}}_{k}\right) f-R(\lambda, \mathcal{L}) f=R\left(\lambda, \tilde{\mathcal{A}}_{k}\right)\left(\mathcal{L}-\tilde{\mathcal{A}}_{k}\right) R(\lambda, \mathcal{L}) f
$$

we deduce

$$
\left\|R\left(\lambda, \tilde{\mathcal{A}}_{k}\right) f-R(\lambda, \mathcal{L}) f\right\|_{p} \leq\left\|R\left(\lambda, \tilde{\mathcal{A}}_{k}\right)\right\|\left\|\mathcal{L} u-\tilde{\mathcal{A}}_{k} u\right\|_{p} \leq\|R(\lambda, \tilde{\mathcal{A}})\|\left\|\mathcal{L} u-\tilde{\mathcal{A}}_{k} u\right\|_{p}
$$

and the claim follows.
Now, let $\lambda_{0} \in \sigma(\mathcal{L})$. Suppose that $\lambda_{0} \in \rho(\tilde{\mathcal{A}})$. Then there exists a radius $r>0$ such that $\lambda \in \rho(\tilde{\mathcal{A}})=\rho\left(\tilde{\mathcal{A}}_{k}\right)$ whenever $\left|\lambda-\lambda_{0}\right|<r$. Take $\lambda$ with $\left|\lambda-\lambda_{0}\right|<r$ and $\operatorname{Re} \lambda>\operatorname{Re} \lambda_{0}$. Then $\lambda$ also belongs to $\rho(\mathcal{L})$ by Theorem 2.2. The formula (4.12) thus yields

$$
\|R(\lambda, \mathcal{L}) f\|_{p} \leq \liminf _{k \rightarrow \infty}\left\|R\left(\lambda, \tilde{\mathcal{A}}_{k}\right) f\right\|_{p} \leq\|R(\lambda, \tilde{\mathcal{A}})\|\|f\|_{p}
$$

for every $f \in L^{p}$. In the limit $\lambda \rightarrow \lambda_{0}$ the left-hand side blows up, whereas the right-hand side remains bounded. By this contradiction, $\lambda_{0}$ belongs to $\sigma(\tilde{\mathcal{A}})$ and consequently $\lambda_{0}-s^{2}\langle Q \xi, \xi\rangle$ to $\sigma(\mathcal{A})$, see (4.10). As $s \in \mathbb{R}$ is arbitrary, we conclude that $\sigma(\mathcal{L})+(-\infty, 0] \subseteq \sigma(\mathcal{A})$.

We next treat the remaining case.
Proposition 4.4. Let $\beta \neq 0$ and $i \beta \in \sigma(B)$, then $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right) \subseteq \sigma\left(\mathcal{A}_{p}\right)$.
In order to show this proposition, we proceed as in Case 2 of the proof of Proposition 4.2 obtaining the same limit operator $\mathcal{C}=\mathcal{A}^{\text {lim }}+\mathcal{L}^{\text {lim }}$, see (4.7). But now we cannot use Theorem 2.9 to determine the spectrum of $\mathcal{A}^{\lim }$, since the $2 \times 2$ drift matrix of $\mathcal{A}^{\text {lim }}$ has the purely imaginary eigenvalues $\pm i \beta$. Instead we directly compute the spectrum of $\mathcal{A}^{\text {lim }}$. We start with a first-order operator that will appear in a scaling limit.

Lemma 4.5. Let $b, s, \mu_{1}, \mu_{2} \in \mathbb{R}$ and set $\mathcal{T}_{\infty} u=b x_{2} D_{1} u-b x_{1} D_{2} u-s^{2}\left(\mu_{1} \frac{x_{1}^{2}}{|x|^{2}}+\mu_{2} \frac{x_{2}^{2}}{|x|^{2}}\right) u$. Let $\mathcal{T}_{\infty, p}$ be the realization of $\mathcal{T}_{\infty}$ in $L^{p}\left(\mathbb{R}^{2}\right)$ endowed with domain $D\left(\mathcal{T}_{\infty}\right)=\left\{u \in L^{p}\left(\mathbb{R}^{2}\right) \mid \mathcal{T}_{\infty} u \in\right.$ $\left.L^{p}\left(\mathbb{R}^{2}\right)\right\}$, where $\mathcal{T}_{\infty} u$ is understood in the sense of distributions. Then, for every $m \in \mathbb{Z}$, the number imb $-s^{2}\left(\mu_{1}+\mu_{2}\right) / 2$ is an eigenvalue of $\mathcal{T}_{\infty, p}$ possessing an eigenfunction $u$ in $C_{c}^{\infty}\left(\mathbb{R}^{2} \backslash\{0\}\right)$.

Proof. In polar coordinates $(\rho, \theta)$, our operator is expressed by

$$
\mathcal{T}_{\infty} u=-b \partial_{\theta} u-s^{2}\left(\mu_{1} \cos ^{2} \theta+\mu_{2} \sin ^{2} \theta\right) u
$$

Let $\varphi \in C_{c}^{\infty}(0, \infty)$ and $m \in \mathbb{Z}$. Set $u(x)=\varphi(|x|) e^{i m \theta} e^{-s^{2}\left(\mu_{1}-\mu_{2}\right) \sin (2 \theta) /(4 b)}$ for $x \in \mathbb{R}^{2}$. Then $u$ belongs to $C_{c}^{\infty}\left(\mathbb{R}^{2} \backslash\{0\}\right)$ and straightforward computations show that $\mathcal{T}_{\infty} u=\lambda u$ with $\lambda=i m b-s^{2}\left(\mu_{1}+\mu_{2}\right) / 2$.

Lemma 4.6. Let $\mathcal{A}^{\diamond}$ be a hypoelliptic Ornstein-Uhlenbeck operator on $\mathbb{R}^{2}$ whose drift matrix $B^{\diamond}$ has the eigenvalues $\pm i \beta$ for $\beta \in \mathbb{R} \backslash\{0\}$. Then $(-\infty, 0]+i \beta \mathbb{Z}=\sigma_{a p}\left(\mathcal{A}_{p}^{\diamond}\right)=\sigma\left(\mathcal{A}_{p}^{\diamond}\right)$.

Proof. We divide the proof in four steps.

1) Put $\mathcal{A}^{\diamond}$ in a canonical form. Let $\mu_{1}$ and $\mu_{2}$ be the two nonnegative eigenvalues of the diffusion matrix $Q^{\diamond}$ of $A^{\diamond}$. There is an invertible matrix $M_{1} \in \mathbb{R}^{2 \times 2}$ such that

$$
M_{1} B^{\diamond} M_{1}^{-1}=\left(\begin{array}{cc}
0 & \beta \\
-\beta & 0
\end{array}\right)
$$

Then take an orthogonal $U_{2} \in \mathbb{R}^{2 \times 2}$ such that $U_{2}\left(M_{1} Q^{\diamond} M_{1}^{T}\right) U_{2}^{T}=D$ for the diagonal matrix $D$ with diagonal elements $\mu_{1}$ and $\mu_{2}$. Since $U_{2}$ is $2 \times 2$ orthogonal, we obtain

$$
U_{2}\left(\begin{array}{cc}
0 & \beta \\
-\beta & 0
\end{array}\right) U_{2}^{T}=\left(\begin{array}{cc}
0 & b \\
-b & 0
\end{array}\right)=: B^{\circ}
$$

where $b= \pm \beta$. The change of variables (3.3) with $M=U_{2} M_{1}$ thus yields

$$
\mathcal{A}^{\circ}=\Phi_{M} \mathcal{A}^{\diamond} \Phi_{M}^{-1}=\mu_{1} D_{11} u+\mu_{2} D_{22} u+b x_{2} D_{1} u-b x_{1} D_{2} u
$$

with $D\left(\mathcal{A}_{p}^{\circ}\right)=\Phi_{M} D\left(\mathcal{A}_{p}^{\diamond}\right)$.
We observe that there are two possible cases: either $\mu_{1}$ and $\mu_{2}$ are both positive, or one of them is positive and the other one zero. In the first case $\mathcal{A}^{\circ}$ is a nondegenerate Ornstein-Uhlenbeck operator, in the second one it is a degenerate hypoelliptic operator.
2) Scale $\mathcal{A}^{\circ}$ by the isometry (4.8). We now set $\phi(x)=|x|$ for all $x=\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}$ in (4.8). Observe that $D \phi(x)=\frac{1}{|x|} x$ and $\left\langle B^{\circ} x, D \phi\right\rangle=0$. For $u \in C_{c}^{\infty}\left(\mathbb{R}^{2} \backslash\{0\}\right)$ and $s \in \mathbb{R}$ the
formula (4.9) thus yields

$$
\begin{aligned}
S^{-1} \mathcal{A}^{\circ} S u(x)= & \mathcal{A}^{\circ} u-s^{2}\left(\mu_{1} \frac{x_{1}^{2}}{|x|^{2}}+\mu_{2} \frac{x_{2}^{2}}{|x|^{2}}\right) u+2 i s\left(\mu_{1} \frac{x_{1}}{|x|} D_{1} u+\mu_{2} \frac{x_{2}}{|x|} D_{2} u\right) \\
& +i s\left(\mu_{1} \frac{x_{2}^{2}}{|x|^{3}}+\mu_{2} \frac{x_{1}^{2}}{|x|^{3}}\right) u=: \mathcal{T} u .
\end{aligned}
$$

By similarity, it is enough to treat $\mathcal{T}$.
3) Show that $i \beta \mathbb{Z}+(-\infty, 0] \subseteq \sigma\left(\mathcal{T}_{p}\right)=\sigma\left(\mathcal{A}_{p}^{\diamond}\right)$. We scale the operator $\mathcal{T}$ through the isometries $V_{k}$ in (4.11), obtaining

$$
\begin{aligned}
V_{k}^{-1} \mathcal{T} V_{k} u= & \frac{1}{k^{2}} \mu_{1} D_{11} u+\frac{1}{k^{2}} \mu_{2} D_{22} u+b x_{2} D_{1} u-b x_{1} D_{2} u-s^{2}\left(\mu_{1} \frac{x_{1}^{2}}{|x|^{2}}+\mu_{2} \frac{x_{2}^{2}}{|x|^{2}}\right) u \\
& +\frac{1}{k} 2 i s\left(\mu_{1} \frac{x_{1}}{|x|} D_{1} u+\mu_{2} \frac{x_{2}}{|x|} D_{2} u\right)+\frac{1}{k} i s\left(\mu_{1} \frac{x_{2}^{2}}{|x|^{3}}+\mu_{2} \frac{x_{1}^{2}}{|x|^{3}}\right) u
\end{aligned}
$$

Set $\mathcal{T}_{k}=V_{k}^{-1} \mathcal{T} V_{k}$. With the limit operator $\mathcal{T}_{\infty}$ from Lemma 4.5, it follows

$$
\mathcal{T}_{k} u \longrightarrow \mathcal{T}_{\infty} u=b x_{2} D_{1} u-b x_{1} D_{2} u-s^{2}\left(\mu_{1} \frac{x_{1}^{2}}{|x|^{2}}+\mu_{2} \frac{x_{2}^{2}}{|x|^{2}}\right) u
$$

in $L^{p}$ as $k \rightarrow \infty$, for every $u \in C_{c}^{\infty}\left(\mathbb{R}^{2} \backslash\{0\}\right)$.
We now argue as in the proof of Proposition 4.2. Take $\lambda_{0} \in \rho\left(\mathcal{T}_{p}\right)$. By similarity, we have $\lambda_{0} \in \rho\left(\mathcal{T}_{k}\right)$ and $\left\|R\left(\lambda_{0}, \mathcal{T}_{k}\right)\right\| \leq C$ for every $k \in \mathbb{N}$. For $u \in C_{c}^{\infty}\left(\mathbb{R}^{2} \backslash\{0\}\right)$ we derive $\|u\|_{p} \leq C\left\|\left(\lambda_{0}-\mathcal{T}_{k}\right) u\right\|_{p}$, and thus $\|u\|_{p} \leq C\left\|\left(\lambda_{0}-\mathcal{T}_{\infty}\right) u\right\|_{p}$ letting $k \rightarrow+\infty$. Recalling that $b= \pm \beta$, Lemma 4.5 implies that $\lambda_{0} \notin i \beta \mathbb{Z}-s^{2}\left(\mu_{1}+\mu_{2}\right) / 2$ for all $s \in \mathbb{R}$. We have this shown the inclusion $i \beta \mathbb{Z}+(-\infty, 0] \subseteq \sigma\left(\mathcal{T}_{p}\right)=\sigma\left(\mathcal{A}_{p}^{\diamond}\right)$.
4) Compute the spectrum of $\mathcal{A}_{p}^{\diamond}$. Theorem $2.2(\mathrm{~d})$ shows that $i \beta \mathbb{Z}$ is the spectrum of the drift operator $\mathcal{L}^{\diamond}=\left\langle B^{\diamond} x, D\right\rangle$. From Proposition 2.8 we deduce $\sigma\left(\mathcal{A}_{p}^{\diamond}\right) \subseteq(-\infty, 0]+i \beta \mathbb{Z}$ and hence $\sigma\left(\mathcal{A}_{p}^{\diamond}\right)=i \beta \mathbb{Z}+(-\infty, 0]$ by step 3$)$. In particular, $\sigma\left(\mathcal{A}_{p}^{\diamond}\right)$ coincides with its topological boundary so that $\sigma\left(\mathcal{A}_{p}^{\diamond}\right)=\sigma_{a p}\left(\mathcal{A}_{p}^{\diamond}\right)$.
Proof of Proposition 4.4. As already pointed out, we proceed as in Case 2 of the proof of Proposition 4.2. The operator $\mathcal{A}$ thus has the form $\mathcal{A}=\operatorname{Tr}\left(Q_{0} D^{2}\right)+\langle T x, D\rangle$ with $T$ as in (3.4). As after (4.6), the functions $J_{r}^{-1} \mathcal{A} J_{r} u$ tend to $\mathcal{C} u$ in $L^{p}$ as $r \rightarrow+\infty$ for $u \in C_{c}^{\infty}$, where $\mathcal{C}=\mathcal{A}^{\text {lim }}+\mathcal{L}^{\lim }$ is defined in (4.7) and it is split into the same operators $\mathcal{A}^{\text {lim }}$ and $\mathcal{L}^{\text {lim }}$. In particular, $\mathcal{A}^{\text {lim }}$ is hypoelliptic. Lemma 4.6 yields the spectral identity $\sigma\left(\mathcal{A}_{p}^{\lim }\right)=\sigma_{a p}\left(\mathcal{A}_{p}^{\lim }\right)=(-\infty, 0]+i \beta \mathbb{Z}$.

We next want to show the equality $\sigma\left(\mathcal{L}_{p}^{\lim }\right)+i \beta \mathbb{Z}=\sigma\left(\mathcal{L}_{p}\right)$. Let $\bar{B}$ denote the coefficient matrix of $\mathcal{L}^{\text {lim }}$. Observe that it is diagonalizable since it has $n-2$ eigenvalues (counted with multiplicities) and that $\sigma(B)=\sigma(\bar{B}) \cup\{ \pm i \beta\}$. Hence, case (b) of Theorem 2.2 does not occur for $\mathcal{L}^{\text {lim }}$. Let $\mathcal{L}^{\lim }$ fall under cases (a) or (c) of Theorem 2.2 (so that $\mathcal{L}$ cannot fall under case (d)). Theorem 2.2 then leads to

$$
\sigma\left(\mathcal{L}_{p}^{\lim }\right)=-\operatorname{tr}(\bar{B}) / p+i \mathbb{R}=-\operatorname{tr}(B) / p+i \mathbb{R}
$$

The asserted equality thus follows from Theorem 2.2. In case (d), Theorem 2.2 and Proposition 2.3 yield

$$
\sigma\left(\mathcal{L}_{p}^{\lim }\right)=\left\{i\left(n_{1} \sigma_{1}+\cdots+n_{m-1} \sigma_{m-1}\right) \mid\left(n_{1}, \cdots, n_{m-1}\right) \in \mathbb{Z}^{m-1}\right\}
$$

If $\beta / \sigma_{1}$ is rational, we infer from these results

$$
\begin{aligned}
\sigma\left(\mathcal{L}_{p}\right) & =\left\{i\left(n_{1} \sigma_{1}+\cdots+n_{m} \beta\right) \mid\left(n_{1}, \cdots, n_{m}\right) \in \mathbb{Z}^{m}\right\} \\
& =\left\{i\left(n_{1} \sigma_{1}+\cdots+n_{m-1} \sigma_{m-1}\right) \mid\left(n_{1}, \cdots, n_{m-1}\right) \in \mathbb{Z}^{m-1}\right\}+i \beta \mathbb{Z} \\
& =\sigma\left(\mathcal{L}^{\lim }\right)+i \beta \mathbb{Z}
\end{aligned}
$$

Otherwise, it follows $\sigma\left(\mathcal{L}_{p}\right)=i \mathbb{R}=\sigma\left(\mathcal{L}^{\text {lim }}\right)+i \beta \mathbb{R}$ as well.
Observe that also $\sigma\left(\mathcal{L}_{p}^{\lim }\right)=\sigma_{a p}\left(\mathcal{L}_{p}^{\lim }\right)$. Take $\lambda_{1} \in \sigma\left(\mathcal{A}_{p}^{\lim }\right)$ and $\lambda_{2} \in \sigma\left(\mathcal{L}_{p}^{\lim }\right)$. As in (4.4), we check that $\mu=\lambda_{1}+\lambda_{2}$ is an approximate eigenvalue for $\mathcal{C}_{p}$; i.e., $(-\infty, 0]+i \beta \mathbb{Z}+\sigma\left(\mathcal{L}^{\lim }\right)=$ $(-\infty, 0]+\sigma\left(\mathcal{L}_{p}\right)$ is contained in $\sigma_{a p}\left(\mathcal{C}_{p}\right)$. Arguing as in (4.5), we finally see that $\sigma_{a p}\left(\mathcal{C}_{p}\right)$ is a subset of $\sigma\left(\mathcal{A}_{p}\right)$.

Propositions 4.2, 4.3 and 4.4 conclude the proof of Theorem 4.1.

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