

On hypoellipticity of degenerate operators in testing and detection problems

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9th Eastern Conference on Mathematical Finance

October 25, 2025

Outline

- 1 Introduction and motivation
- 2 Problem formulation
- 3 Main Results
- 4 Examples

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1 Introduction and motivation

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The classic 1D sequential testing

- A stopper observes a 1D Brownian motion with drift

$$dX_t = \theta t + dW_t, \quad X_0 = 0.$$

with $P(\theta = 1) = 1 - \theta = 0 = \pi \in [0, 1]$.

- Want to test hypotheses for its drift, e.g., $H_0 : \theta = 0$ vs $H_1 : \theta = 1$.
- We are penalized for making a mistake, and have a constant observation cost c per unit time.
- “Sequential testing problem”, can be formulated as an optimal stopping problem in the posterior probability process, as

$$V(\pi) = \inf_{\tau} \mathbb{E}[\Pi_{\tau} \wedge (1 - \Pi_{\tau}) + c\tau]$$

where $\Pi_t := \mathbb{P}(\theta = 1 | \mathcal{F}_t^X)$.

- Can be formulated as a free-boundary problem and solved explicitly [cf. Shiryaev (1978)].

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A challenge when in higher dimensions: degeneracy

Multi-dimension setting? many possibilities. Let's for example, look at the following problem:

- Consider now the Brownian motion has 3 possible drifts instead of 2:
 $P(\theta = i) = \pi_i, i \in \{0, 1, 2\}.$
- The sufficient statistics in this case is (Π^1, Π^2) .
- There are two underlying coordinates but only one underlying Brownian source. The operator is **degenerate elliptic**.
- Π_t^i are functions of (t, X_t) . Can formulate it in the (t, x) –coordinate: uniformly parabolic [cf. [Zhitlukhin and Shiryaev\(2011\)](#)].

A challenge when in higher dimensions: degeneracy

What about the following cases when they are degenerate in π ?

- When θ can change its value at exponential times. e.g. classic quickest detection:

$$dX_t = 1_{\theta \geq t} dt + dW_t.$$

The posterior process $\Pi_t := \mathbb{P}(\theta > t | \theta > 0)$ depends on the whole path. No longer possible to formulate in (t, x) .

- When the problem is X -dependent. e.g., a “hiring problem” application:

$$X_t = \theta t + W_t,$$

and the payoff upon stopping at τ being

$$e^{-r\tau} X_\tau.$$

After filtering, the problem has (X, Π) as its state.

These motivate us to study properties of these degenerate cases.

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Problem formulation

The ingredients of our problem:

- A continuous time Markov chain $(\theta_t)_{t \geq 0}$ taking values in $\{0, 1, \dots, n\}$ with generator $Q = (q_{i,j})_{i,j \in \bar{n}}$,
- k -dimensional Brownian motion $W = (W^1, \dots, W^k)$ independent of θ .
- $k < n$.

We consider

$$dX_t = \sum_{j=0}^n 1_{\theta_t=j} \lambda_j dt + dW_t, \quad X_0 = 0. \quad (1)$$

and problems of the form

$$V = \sup_{\tau \in \mathcal{T}} \mathbb{E}_{\pi, x} \left[e^{- \int_0^\tau r(\Pi_s, X_s) ds} g(\Pi_\tau, X_\tau) + \int_0^\tau e^{- \int_0^t r(\Pi_s, X_s) ds} h(\Pi_t, X_t) dt \right]. \quad (2)$$

- Here g, h, r are continuous, $r \geq 0$, $\lambda_i \in \mathbb{R}^k$, $i \in \{0, 1, \dots, n\}$.
- The posterior Pi lives on the n -dimensional simplex P_{n+1}

$$\Pi_t^i = \mathbb{P}_\pi(\theta_t = i \mid \mathcal{F}_t^X) \quad \text{for } i \in \{0, \dots, n\}$$

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Problem formulation

- By standard filtering theory, Π_t has dynamics:

$$d\Pi_t^j = \underbrace{\sum_{i=0}^n q_{ij} \Pi_t^i dt}_{\text{Drift from } Q} + \underbrace{\Pi_t^j (\lambda_j - \bar{\lambda}_t) \cdot d\tilde{W}_t}_{\text{Diffusion from } W_t}$$

- where $\bar{\lambda}_t = \sum_{i=0}^n \lambda_i \Pi_t^i$, and \tilde{W}_t is the innovation process.
- The stopping problem is governed by the infinitesimal generator \mathcal{L}_π for this Π_t process

$$\mathcal{L}_\pi = \underbrace{\frac{1}{2} \sum_{i,j=0}^n \pi_i \pi_j (\lambda_i - \bar{\lambda}) \cdot (\lambda_j - \bar{\lambda}) \frac{\partial^2}{\partial \pi_i \partial \pi_j}}_{\text{Diffusion (degenerate)}} + \underbrace{\sum_{i,j=0}^n q_{ij} \pi_i \frac{\partial}{\partial \pi_j}}_{\text{Drift (from } Q\text{)}}$$

- **It degenerates everywhere.**

Hypoellipticity and the Hörmander's condition

- **Question:** Even if the operator is not elliptic, can we still recover regularity (e.g., smoothness) of the value function?
- **Our hope:** Hypoellipticity (the property that u smooth if $\mathcal{L}u$ smooth).
- **Intuition:** The operator may be degenerate, but the randomness "spreads" through the system.
- The "missing directions" from the k -dimensional diffusion might be restored via iterated **Lie brackets**.

Hörmander (1967)

Write $\mathcal{L}_\pi = \sum_{r=1}^k D_r^2 + D_0$, where D_i 's are C^∞ vector fields. If

$$Lie(D_0, D_1, \dots, D_K)$$

spans the tangent space of the simplex at every point in $int(P_{n+1})$, the Hörmander's condition is satisfied, and the operator \mathcal{L}_π is hypoelliptic.

Existing work

- [cf. Caffarelli and Friedman (1981)] studied the problem where $n = k$, with

$$g(\pi) = a_0(1 - \pi_0) \wedge \cdots \wedge a_n(1 - \pi_n) \quad h(\pi) = \sum_{i=0}^n c_i \pi_i.$$

- They commented on the case where $k < n$ and gave the 1D, 3 drift example.
- Few literature in the filtering field: [cf. Peskir (2022), Ernst et al (2022)]

Our goal: characterize when the Hörmander's condition holds for \mathcal{L}_π .

A better coordinate: Φ_t

- We do a change of coordinate to the **posterior likelihood process** Φ_t :

$$\Phi_t^i = \frac{\Pi_t^i}{\Pi_t^0} \quad \text{for } i = 1, \dots, n \quad (\text{Note: } \Phi_t^0 \equiv 1)$$

- This map is a C^∞ -diffeomorphism from $\text{int}(P_{n+1}) \rightarrow (0, \infty)^n$.
Hypoellipticity is preserved.
- Define $a_i \in \mathbb{R}^k$ for $i = 1, \dots, n$: $a_i := \lambda_i - \lambda_0$, and $\Sigma_{ij} = a_i \cdot a_j$.

Denote by $y(\phi) = \sum_{i=0}^n \phi_i$, the generator \mathcal{L} for the Φ process is:

$$\begin{aligned} \mathcal{L} = & \frac{1}{2} \sum_{i,j=1}^n \Sigma_{ij} \phi_i \phi_j \frac{\partial^2}{\partial \phi_i \partial \phi_j} + \frac{1}{y(\phi)} \sum_{i,j=1}^n \Sigma_{ij} \phi_i \phi_j \frac{\partial}{\partial \phi_j} \\ & + \sum_{j=1}^n \sum_{i=0}^n (q_{ij} - q_{i0} \phi_j) \phi_i \frac{\partial}{\partial \phi_j} \end{aligned}$$

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Main Results: characterization in the testing case (Q=0)

We first write \mathcal{L} in the "sum of squares" form $\mathcal{L} = D_0 + \frac{1}{2} \sum_{r=1}^k D_r^2$.

- **Key observation:**

- ① The diffusion fields commute:

$$[D_r, D_u] = 0 \quad \text{for all } r, u \in \{1, \dots, k\}$$

- ② The bracket of the drift and a diffusion field **stays in the diffusion span**:

$$[D_0, D_u] = \sum_{s=1}^k c_s(\phi) D_s \in \text{span}\{D_1, \dots, D_k\}$$

Theorem 0

Let $Q = 0$. Let $A = (a_1, \dots, a_n) \in \mathbb{R}^{k \times n}$.

If $n > k + 1$, the Hömander's condition **FAILS**.

If $n = k + 1$, the Hömander's condition **HOLDS** if and only if $\text{rank}(A) = k$ and the vector $(\|a_1\|^2, \dots, \|a_n\|^2)$ is not in the rowspace of A .

- Define "switch" field $J := \sum_{j=1}^n \sum_{i=0}^n (q_{i,j} - q_{i,0} \phi_j) \phi_i \partial_{\phi_j}$
- **The Key:** The Lie bracket of J with the diffusion fields D_r creates *new* vector fields.

$$[J, D_r], \quad [D_s, [J, D_r]], \quad \text{etc.}$$

Theorem 1 (Sufficient Condition 1)

The Hömander's condition holds if: (1) The drift-difference vectors a_1, \dots, a_n are pairwise distinct, and (2) **For each** coordinate $i \in \{1, \dots, n\}$, there exists some state $m \neq i$ such that $q_{mi} > 0$.

Remark: This is much weaker than Q being irreducible.

Theorem 2 (Sufficient Condition 2)

Assume $q_{mj} = 0$ for all $m \neq j$ and that there is at least one j such that $q_{j0} > 0$ (at least one state can jump back to 0).

Define the $n \times (k + 2)$ augmented matrix \tilde{A} :

$$\tilde{A} = \begin{bmatrix} a_{11} & \dots & a_{k1} & \|a_1\|^2 & 1 \\ \vdots & \ddots & \vdots & \vdots & \vdots \\ a_{1n} & \dots & a_{kn} & \|a_n\|^2 & 1 \end{bmatrix}$$

Then, $\dim \text{Lie}(D_0^J, \dots, D_k) = \min(\text{rank}(\tilde{A}), n)$. Hypoellipticity holds if $\text{rank}(\tilde{A}) = n$, which can **never** hold if $n > k + 2$.

Consequences of Hörmander

Parabolic Hörmander's condition is immediate:

- Define $\bar{D}_0 := -\partial_t + D_0$, $\dim \text{Lie}\{D_0^J, D_1, D_2, \dots, D_k\} = n$, then $\dim \text{Lie}\{\bar{D}_0^J, D_1, D_2, \dots, D_k\} = n+1$.
- It follows that the process $(\Phi_t)_{t \geq 0}$ is strong Feller.
- This allows us to deal with t –dependent problems.

Hörmander holds in the (ϕ, x) coordinate:

- The operator $\mathcal{L}_{\phi, x}$ for (ϕ, x) –dependent problems is hypoelliptic on $(0, \infty)^n \times \mathbb{R}^k$ if and only if the operator \mathcal{L} is hypoelliptic on $(0, \infty)^n$.
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Value function regularity in \mathcal{C}

Let C be the **continuation region** (where $V > g$).

- If the Hömander condition holds, and the running payoff $h \in C^\infty(C)$, then the value function V is also $C^\infty(C)$.
- If $r, h \in C^{0,\alpha}(C)$ for $\alpha \in (0, 1)$, then the value function $V \in C^{2,\alpha}(C)$.

But no "smooth fit" implied:

- All these regularity results (C^∞ or $C^{2,\alpha}$) hold **only** in the \mathcal{C} .
- In particular, hypoellipticity alone **does not imply** the "smooth fit" condition.

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Example 1: Detecting the change time in multiple coordinate

- We observe an N -dim BM. At time η , K out of N coordinates gain a drift μ .
- There are N Brownian coordinates and $\binom{N}{K}$ possible drifts.
- The generator Q has a specific structure: only the first row is non-zero.

$$Q = \begin{bmatrix} -\binom{N}{K}\lambda & \lambda & \dots & \lambda \\ 0 & 0 & \dots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \dots & 0 \end{bmatrix}$$

- The drift-difference vectors a_j are all distinct. ✓
- All $q_{0i} > 0$ ✓.
- The operator is **hypoelliptic** by Theorem 1.

[cf. Ernst et al (2022)]

Example 2: Signal tracking with regime switching

- We observe $X_t = \sum_{j=0}^n 1_{\theta_t=j} \lambda_j t + W_t$. with

$$Q = \begin{bmatrix} -\sum_{i=1}^n q_i & q_1 & q_2 & \dots & q_n \\ p_1 & -p_1 & 0 & \dots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ p_k & 0 & 0 & \dots & -p_k \end{bmatrix} \quad (3)$$

with $p_i, q_j > 0$ for all $i \in \{1, k\}, j \in \{1, n\}$.

- Application-wise: monitoring signals from a radar with different levels of disorder. We can consider, e.g.,

$$\inf_d \mathbb{E} \left[\int_0^T e^{-rt} c(d_t, \theta_t) dt \right]$$

with

$$c(d, \theta) = 1_{d=\theta} \sum_{i=1}^n c^1(\theta) + 1_{d \neq \theta} \sum_{i=1}^n c^2(\theta).$$

- The operator is **hypoelliptic** by Theorem 1 (same as Example 1).

Example 3: Byzantine detection

A detection problem with possible corrupted sensor:

- We observe two processes X^1, X^2 :

$$X_t^1 = \mu_0 1_{\eta < t} + \mu_1 1_{\eta \leq t} + W_t^1,$$
$$X_t^2 = m_0 1_{\eta < t} + m_1 1_{\eta \leq t} + W_t^2,$$

with $\mathbb{P}_0(\eta = 0) = p$, $\mathbb{P}_0(\eta > t) = (1 - p)e^{-rt}$, $t \geq 0$.

- The unknown state θ has 4 possibilities:

$$\theta_0 = \begin{cases} (\mu_1, m_1), & \text{both channels are affected at } t = 0 \\ (\mu_1, m_0), & X^1 \text{ is affected at time } t = 0, \\ (\mu_0, m_1), & X^2 \text{ is affected at time } t = 0, \\ (\mu_0, m_0), & \text{no channels are affected at time } t = 0, \end{cases} \quad (4)$$

We call it Byzantine because of the “Byzantine general problem”.

Example 3: Byzantine detection

- The drift matrix

$$A = - \begin{bmatrix} \mu_1 - \mu_0 & 0 & \mu_1 - \mu_0 \\ 0 & m_1 - m_0 & m_1 - m_0 \end{bmatrix}$$

has full rank (2), but the vector $(\|a_1\|^2, \|a_2\|^2, \|a_3\|^2)$ is always in the row-space of A .

- The generator matrix Q as

$$Q = \begin{bmatrix} 0 & 0 & 0 & 0 \\ \lambda & -\lambda & 0 & 0 \\ \lambda & 0 & -\lambda & 0 \\ \lambda & 0 & 0 & -\lambda \end{bmatrix}.$$

Theorem 3 applies, the augmented matrix

$$\tilde{A} = - \begin{bmatrix} \mu_1 - \mu_0 & 0 & (\mu_1 - \mu_0)^2 & 1 \\ 0 & m_1 - m_0 & (m_1 - m_0)^2 & 1 \\ \mu_1 - \mu_0 & m_1 - m_0 & (\mu_1 - \mu_0)^2 + (m_1 - m_0)^2 & 1 \end{bmatrix}$$

has rank 3, Hörmander condition holds. In working progress.

Thank you!