Nonexistence of finite dimensional filters for conditional statistics of the cubic sensor problem
by
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ABSTRACT

Consider the cubic sensor $d x=d w, d y=x^{3} d t+d v$ where $w, v$ are two independent brownian motions. Given a function $\phi(x)$ of the state $x$ let $\phi_{t}(x)$ denote the conditional expection given the observations $y_{s}, 0 \leq s \leq t$. This paper consists of a rather detailed discussion and outline of proof of the theorem that for nonconstant $\phi$ there can not exist a recursive finite dimensional filter for $\hat{\phi}$ driven by the observations.

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KEY WORDS \& PHRASES: cubic sensor, recursive filter, robust filtering, Weyt Lie algebra
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1. INTRODUCTION

The cubic sensor problem is the problem of determining conditional statistics of the state of a one dimensional stochastic process $\left\{x_{t}: t \geq 0\right\}$ satisfying
(1.1) $d x=d w, \quad x_{0}=x^{i n}$
with w a Wiener process, independent of $x^{\text {in }}$, given the observation process $\left\{y_{t}: t \geq 0\right\}$ satisfying
(1.2) $\quad d y=x^{3} d t+d v, \quad y_{0}=0$
where v is another Wiener process independent of w and $\mathrm{x}^{\text {in }}$. Given a smooth function $\phi: \mathbb{R} \rightarrow \mathbb{R}$ let $\oint_{t}$ denote the conditional expection

$$
\begin{equation*}
\Phi_{t}=\phi\left(x_{t}\right)=E\left[\phi\left(x_{t}\right) \mid y_{s}, 0 \leq s \leq t\right] \tag{1.3}
\end{equation*}
$$

By definition a smooth finite dimensional recursive filter for $\phi_{t}$ is a dynamical system on a smooth finite dimensional manifold M governed by an equation

$$
\begin{equation*}
d z=\alpha(z) d t+\beta(z) d y, \quad z_{0}=z^{\text {in }} \tag{1.4}
\end{equation*}
$$

driven by the observation process, together with an output map

$$
\begin{equation*}
\gamma: M \rightarrow \mathbb{R} \tag{1.5}
\end{equation*}
$$

such that, if $z_{t}$ denotes the solution of (1.4),

$$
\begin{equation*}
\gamma\left(z_{t}\right)=\phi_{t} \text { a.s. } \tag{1.6}
\end{equation*}
$$

Roughly speaking one now has the theorem that for nonconstant $\phi$ such filters cannot exist. For a more precise statement of the theorem see 2.10 below.

It is the purpox of this note to give a fairly detailed outline of the proof of this theorem and to discuss the structure of the proof. That is the general principles underlying it. The full precise details of the analytic and realization theoretic parts of the proof will appear in [Sussmann 1983a, 1983b], the details of the algebraic part of the proof can be found in [Hazewinkel-Marcus, 1982]. An alternative much better and shorter proof of the hardest bit of the algebraic part will appear in [Stafford, 1983].
2. SYSTEM THEORETIC PART. I: PRECISE FORMULATION OF THE THEOREM

### 2.1 The setting

The precise system theoretic - probabilistic setting which we shall use for the cubic sensor filtering problem is as follows
(i) $(\Omega, A, P)$ is a probability space
(ii) ( $A_{t}: 0 \leq t$ ) is an increasing family of $\sigma$-algebras
(iii) (w,v) is a two-dimensional standard Wiener process adapted to the $A_{t}$.
(iv) $x=\left\{x_{t}: t \geq 0\right\}$ is a process which satisfies $d x=d w$, i.e.
(2.1) $\quad x_{t}=x_{0}+w_{t} \quad$ a.s. for each $t$
(v) $x_{0}$ is $A_{0}$-measurable and has a finite fourth moment
(vi) $\left\{y_{t}: t \geq 0\right\}$ is a process which satisfies $d y=x^{3} d t+d v$, i.e.

$$
\begin{equation*}
y_{t}=\int_{0}^{t} x_{s}^{3} d s+v_{t} \quad \text { a.s. for each } t \tag{2.2}
\end{equation*}
$$

(vii) the processes $v, w, x$, $y$ all have continuous sample paths, so that in particular (2.1) and (2.2) actually hold and not just almost surely. (More precisely one can always find if necessary modified versions of $v, w, x, y$ such that (vii) (also) holds).

### 2.3. The filtering problem

Let $y_{t}, t \geq 0$ be the $\sigma$-algebra generated by the $y_{s}, 0 \leq s \leq t$ and let $\phi: \mathbb{R} \rightarrow \mathbb{R}$ be a Borel measurable function. Then the filtering problem (for this particular $\phi$ ) consists of determining $E\left[\phi\left(x_{t}\right) \mid Y_{t}\right]$.

### 2.4. Smooth finite dimensional filters

Consider a (Fisk-Stratonovic) stochastic differential equation

$$
\begin{equation*}
d z=\alpha(z) d t+\beta(z) d y, \quad z \in M, \tag{2.5}
\end{equation*}
$$

where $M$ is a finite dimensional smooth manifold and $\alpha$ and $\beta$ are smooth vectorfields on M. Let there also be given an initial state and a smooth output map

$$
\begin{equation*}
z^{\text {in }} \in M, \quad \gamma: M \rightarrow \mathbb{R} \tag{2.6}
\end{equation*}
$$

The equation (2.5) together with the initial condition $z(0)=z^{\text {in }}$ has a solution $z=\left\{z_{t}: t \geq 0\right\}$ defined up to a stopping time $T$, which satisfies

$$
\begin{equation*}
0<T \leq \infty \text { a.s., } \quad\{\omega \mid T(\omega)>t\} \in Y_{t}, \text { for } t \geq 0 \tag{2.7}
\end{equation*}
$$

Moreover there is a unique maximal solution, i.e. one for which the stopping time T is a.s. $\geq \mathrm{T}_{1}$ if $\mathrm{T}_{1}$ is the stopping time of an arbitrary other solution $z_{1}$. In the following $z=\left\{z_{t}: t \geq 0\right\}$ denotes such a maximal solution.

The system given by (2.5), (2.6) is now said to be a smooth finite dimensional filter for the cubic sensor (2.1) (i) - (vii) if for y equal to the observation process (2.2) the solution $z$ of (2.5) satisfies

$$
\begin{equation*}
E\left[\phi\left(x_{t}\right) \mid Y_{t}\right]=\gamma\left(z_{t}\right) \text { a.s. on }\{\omega \mid T(\omega)>t\} . \tag{2.8}
\end{equation*}
$$

### 2.9. Statement of the theorem

With these notions the main theorem of this note can be stated as:
2.10. THEOREM. Consider the cubic sensor 2.1. (i) - (vii); i.e. assume that these conditions hold. Let $\phi: \mathbb{R} \rightarrow \mathbb{R}$ be a Borel measurable function which satisfies for some $\beta \geq 0$ and $0 \leq r<4$

$$
\begin{equation*}
|\phi(x)| \leq \exp \left(\beta|x|^{r}\right), \quad-\infty<x<\infty . \tag{2.11}
\end{equation*}
$$

Assume that $\phi$ is not almost everywhere equal to a constant. Then there exists no smooth finite dimensional filter for the conditional statistic $E\left[\phi\left(x_{t}\right) \mid Y_{t}\right]$.
3. SYSTEM THEORETIC PART. II: THE HOMOMORPHISM PRINCIPLE AND OUTLINE OF THE PROOF (HEURISTICS)

### 3.1. The Duncan-Mortensen-Zakai equation

Consider a nonlinear stochastic dynamical system

$$
\begin{equation*}
d x_{t}=f\left(x_{t}\right) d t+G\left(x_{t}\right) d w_{t}, \quad x^{i n}=x_{0}, \quad x_{t} \in \mathbb{R}^{n}, \quad w_{t} \in \mathbb{R}^{m} \tag{3.2}
\end{equation*}
$$

where $w_{t}$ is a standard Brownian motion independent of the initial random variable $x^{i n}$ and where $f$ and $G$ are appropriate vector valued and matrix valued functions. Let the observations be given by

$$
\begin{equation*}
d y_{t}=h\left(x_{t}\right) d t+d v_{t}, \quad y_{t} \in \mathbb{R}^{p} \tag{3.3}
\end{equation*}
$$

where $\mathrm{v}_{\mathrm{t}}$ is another standard Brownian motion independent of w and $\mathrm{x}^{\text {in }}$. Let $\hat{\mathrm{x}}_{\mathrm{t}}$ denote the conditional expectation

$$
\begin{equation*}
\hat{x}_{t}=E\left[x_{t} \mid Y_{t}\right]=E\left[x_{t} \mid y_{s}, 0 \leq s \leq t\right] \tag{3.4}
\end{equation*}
$$

where $Y_{t}$ is the $\sigma$-algebra generated by the $y_{s}, 0 \leq s \leq t$. Let $p(x, t)$ be the density of $\hat{x}_{t}$ where it is assumed (for the purposes of this heuristic section) that $p(x, t)$ exists and is sufficiently smooth as a function of $x$ and $t$. Then an unnormalized version $\rho(x, t)$ satisfies the Duncan-Mortensen-Zakai equation

$$
\begin{align*}
d \rho(x, t) & =\left(\frac{1}{2} \sum_{i, j} \frac{\partial^{2}}{\partial x_{i} \partial} j_{j}\left(\left(G G^{T}\right)_{i j}-\sum_{i} \frac{\partial}{\partial x_{i}} f_{i}-\frac{1}{2} \sum h_{j}^{2}\right) \rho(x, t) d t\right.  \tag{3.5}\\
& +\sum h_{j} \rho(x, t) d y_{j t}, \quad \rho(x, 0)=\text { density of } x^{i n}
\end{align*}
$$

where $h_{j}=h_{j}(x)$ is the $j$-th component of $h,\left(G G^{T}\right)_{i, j}$ is the (i,j)-th entry of the product of the matrix $G(x)$ with its transpose and $f_{i}=f_{i}(x)$ is the i-th component of $f(x)$. The equation (3.5) is a stochastic partial differential equation in Fisk-Stratonovic form. In the case of the cubic sensor (2.1), (2.2) (or (1.1), (1.2)) the equation becomes

$$
\begin{equation*}
d \rho(x, t)=\left(\frac{1}{2} \frac{d^{2}}{d x^{2}}-\frac{1}{2} x^{6}\right) \rho(x, t) d t+x^{3} \rho(x, t) d y \tag{3.6}
\end{equation*}
$$

### 3.7. The homomorphism principle

Now assume for a given $\phi: \mathbb{R}^{\mathrm{n}} \rightarrow \mathbb{R}$ we have a smooth finite dimensional filter

$$
\begin{equation*}
d z=\alpha(z) d t+\sum_{j} \beta_{j}(z) d y_{j}, \quad z_{0}=z^{i n}, \quad \gamma: \mathbb{R}^{n} \rightarrow \mathbb{R} \tag{3.8}
\end{equation*}
$$

to calculate the statistic $\phi_{t}=E\left[\phi\left(x_{t}\right) \mid Y_{t}\right]$. I.e. $\phi_{t}=\gamma\left(z_{t}\right)$ a.s. if $z_{t}$ is the solution of (3.8). The equation (3.8) is to be interpreted in the Statonovič sense.

Then, very roughly, we have two ways to process an observation path $y^{\omega}: s \rightarrow y_{s}(\omega), 0 \leq s \leq t$ to give the same result. One way is by means of the filter (3.8), the other way is by means of the infinite dimensional system (3.5) (defined on a suitable space of functions) coupled with the output map

$$
\begin{equation*}
\Phi: \psi \rightarrow\left(\int \psi(\mathrm{x}) \mathrm{dx}\right)^{-1} \int \psi(\mathrm{x}) \phi(\mathrm{x}) \mathrm{dx} \tag{3.9}
\end{equation*}
$$

Assuming that (3.8) is observable, deterministic realization theory [Sussmann 1977] then suggest that there exists a smooth map $F$ from the reachable part (from $\rho(x, 0)$ ) of (3.6) to the reachable part of (3.8), which takes the vectorfields of (3.6) to the vectorfields of (3.6) and which is compatible with the output maps $\gamma$ and (3.9). The operators in (3.6) define linear vectorfields in the state space of (3.6) (a space of functions). Let $L_{0}, L_{1}, \ldots, L_{p}$ be the operators occuring in (3.5) so that $\mathrm{d} \rho=\mathrm{L}_{0} \rho \mathrm{dt}+\mathrm{L}_{1} \rho \mathrm{dy}_{1}+\ldots+\mathrm{L}_{\mathrm{p}} \rho \mathrm{dy}_{\mathrm{p}}$. The Lie algebra of differential operators generated by $L_{0}, \ldots, L_{p}$ is called the estimation Lie algebra, and is
denoted $L(\Sigma)$. The idea of studying this Lie algebra to find out things about filtering problems is apparently due to both Brockett and Mitter, cf e.g. [Brockett 1981] and Mitter [1981] and the references in these two papers.

Let $\mathrm{L} \rightarrow \tilde{\mathrm{L}}$ be the map which assigns to an operator the corresponding linear vectorfield (analogous to the map which assigns to an $n \times n$ matrix $A=\left(a_{i j}\right)$ the linear vector field $\sum a_{i j} x_{i} \frac{\partial}{\partial x_{j}}$ as $\left.\mathbb{R}^{n}\right)$. Then $L \leftrightarrow-\tilde{L}$ is a homomorphism of Lie algebras. Further $F$ induces a homomorphism of Lie algebras dF: $\tilde{\mathrm{L}}_{0} \rightarrow \alpha, \tilde{\mathrm{~L}}_{\mathrm{i}} \rightarrow \beta_{i}, i=1, \ldots, p$. Thus the existence of a finite dimensional filter should imply the existence of a homomorphism of Lie algebras $L(\Sigma) \rightarrow V(M)$ where $V(M)$ is the Lie algebra of smooth vectorfields on a smooth finite dimensional manifold $M$. This principle, originally enunciated by Brockett, has come to be called the homomorphism principle.

### 3.10. Pathwise filtering (robustness)

As it stands to remarks in 3.7 above are quite far from a proof of the homomorphism principle. First of all (3.6) and (3.8) are stochastic differential equations and as such they have solutions defined only almost everywhere. The first thing to do to remedy this situation is to show that these equations make sense and have solutions pathwise so that they can be interpreted as processing devices which accept an observation path $y:[0, t] \rightarrow \mathbb{R}^{p}$ and produce outputs $\phi_{t}(y)$ as a result. Another reason for looking for pathwise robust versions which is most important for actual applications, lies in the observation that actual physical observation paths will be piece-wise differentiable and that the space of all such paths is of measure zero in the probability space of paths underlying (3.6) and (3.8). Cf. [Clark, 1978].

Another difficulty in using the remarks of 3.7 above to establish a general homomorphism principle lies in the fact that (3.6) evolves on an infinite dimensional state space. A different approach to the establishing of homomorphism principles (than the one used in this paper) is described in [Hijab, 1982].

### 3.11. On the proof of theorem 2.10

In this paper the following route is followed to establish the homomorphism principle for the case of the cubic sensor. First for suitable $\phi: \mathbb{R} \rightarrow \mathbb{R}$ it is established that there exists a robust pathwise version of the functional $\phi_{t}$. More precisely if $C_{t}$ is the space of continuous functions $[0, t] \rightarrow \mathbb{R}$ then it is shown that there exists a functional $\Delta_{t}^{\phi}: C_{t} \rightarrow \mathbb{R}$ such that (proposition 4.15)

$$
\begin{equation*}
\Phi_{t}=\frac{\Delta_{t}^{\phi}(y)}{\Delta_{t}^{1}(y)} \quad \text { a.s. if } \quad y=y^{\omega} \tag{3.12}
\end{equation*}
$$

The next step is to show that $\Delta_{t}^{\phi}(y), y \in C_{t}$ is given by a density $n_{t}(y)(x)$ so that $\Delta_{t}^{\phi}(y)=\int n_{t}(y)(x) \phi(x) d x$ and to show that $n_{t}(y)(x)$ is smooth (as a function of $x$ ).

The next step is to use that there exist (up to a stopping time) pathwise and robust solutions of stochastic differential equations like (3.8). Robustness of both (3.6) and (3.8) then gives the central equality (4.24) anywhere (not just a.s.), that is:

$$
\begin{equation*}
\frac{\Delta_{t}^{\phi}(y)}{\Delta_{t}^{l}(y)}=\gamma\left(z_{t}(y)\right), \quad y \in C_{t} \tag{3.13}
\end{equation*}
$$

The next step is to prove results about the smoothness properties of the density $n_{t}(y)$ as a function of $t_{1}, \ldots, t_{m}$ for paths $y$ such that $u=\dot{y}$ is of the bang-bang type: $u(s)=\bar{u}_{m} \in \mathbb{R}$ for $0 \leq t<t_{m}$, equal to $\bar{u}_{m-1}$ for $t_{m} \leq t<t_{m}+t_{m-1}$, etc. ... and to observe that $(t, x) \rightarrow n_{t}(y)(x)$ satisfies the $D M Z$ equation (3.6). This permits to write down and calculate the result of applying $\left.\frac{\partial^{m}}{\partial t_{1} \ldots \partial t_{m}} \right\rvert\, t_{1}=\ldots=t_{m}=0$ to both sides of (3.13) and gives a relation of the type

$$
\begin{equation*}
\left(A\left(\bar{u}_{m}\right) \ldots A\left(\bar{u}_{1}\right) \gamma\right)(z)=\tilde{L}\left(\bar{u}_{m}\right) \ldots \tilde{L}\left(\bar{u}_{1}\right) \Phi\left(\psi_{z}\right) \tag{3.14}
\end{equation*}
$$

where $A(\bar{u})$ is the vectorfield $\alpha+\bar{u} \beta, L(\bar{u})$ the operator $L_{0}+\bar{u}_{L}=$ $=\left(\frac{1}{2} \frac{d^{2}}{d x^{2}}-\frac{1}{2} x^{6}\right)+\overline{u x}{ }^{3}$ and $\tilde{L}(\bar{u})$ the linear vectorfield associated to $L(\bar{u})$ $\Phi$ the functional (3.9), and $\psi_{z}$ a function corresponding to $z$, cf. 5.10.

A final realization theoretic argument having to do with reducing the
filter dynamical system (3.8) to an equivalent observable and reachable system then establishes the homomorphism principle in the case of the cubic sensor and the fact that if the homomorphism is zero $\phi$ was a constant.

The remaining algebraic part of the proof consists of two parts:
(i) a calculation of $L(\bar{z})$ for the cubic sensor. It turns out that $L(\bar{z})$ is in this case equal to the Heisenberg-Weyl algebra $W_{1}$ of all differential operators (any order) in $x$ with polynomial coefficients.
(ii) the theorem that if $V(M)$ is the Lie algebra of smooth vectorfields on a smooth finite dimensional manifold and $\alpha: W_{1} \rightarrow V(M)$ a homomorphism of Lie algebras then $\alpha=0$.

## 4. ANALYTIC PART

### 4.1. The space of functions $E$

Let $E$ denote the space of all Borel functions $\phi: \mathbb{R} \rightarrow \mathbb{R}$ such that there exists constants $C \in \mathbb{R}, \alpha \in \mathbb{R}$ and $r, 0 \leq r<4$ such that
(4.2) $|\phi(x)| \leq C \exp \left(\alpha|x|^{r}\right)$ for all $x \in \mathbb{R}$.

The space $E(\alpha, r)$ is the normed space of all Borel functions for which
(4.3) $\quad\|\phi\|_{(\alpha, r)}=\sup \left\{|\phi(x)| \exp \left(-\alpha|x|^{r}\right) ; x \in \mathbb{R}\right\}$
is finite. The space $E$ is the union of the $E(\alpha, r)$ and is topologized as such, i.e. as the inductive limit of the $E(\alpha, r)$.

### 4.2. Some bounds

Let $C_{t}$ be the space of continuous functions $[0, t] \rightarrow \mathbb{R}$ such that $y(0)=0$ and $C_{t}^{1}$ the functions of class $C^{1}$ in $C_{t}$, and $C_{t}^{l \#}$ the functions $y \in C_{t}$ which are piecewise $C^{1}$. Let $H$ be the space of functions on $\Omega$
(4.5) $\quad H=1 \leq \bigcap_{\mathrm{p}}<\infty \quad \mathrm{L}^{\mathrm{P}}(\Omega, \mathrm{A}, \mathrm{P})$.

For the processes $x$ and $w$ of $2.1, \beta, r \in \mathbb{R}$, and a given $y \in C_{t}$ we define

$$
\begin{equation*}
U(\beta, r, t, y)=\exp \left(\beta\left|x_{t}\right|^{r}+y(t) x_{t}^{3}-3 \int_{0}^{t} y(s) x_{s}^{2} d w_{s}-3 \int_{0}^{t} y(s) x_{s} d s-\frac{1}{2} \int_{0}^{t} x_{s}^{6} d s\right) \tag{4.6}
\end{equation*}
$$

The reason for considering this expression becomes clearer below. The $\mathrm{U}(0,0, \mathrm{t}, \mathrm{y})$ occur in a slightly modified version of the Kallianpur-Stiebel formula for $\hat{\phi}_{t}$. The very slight modification (an integration by parts to remove $d_{s}$ ) sees to it that the formula makes pathwise sense for continuous sample paths $y_{s}(\omega)$. The number 4 of 4.1 above gets into the picture as a result of wanting to keep $\left.E\left[\phi\left(x_{t}\right) U \mid 0,0, t, y\right)\right]$ bounded for bounded $y$ and $t$.
4.7. Lemma. $U(\beta, r, t, y) \in H$ if $0 \leq \beta, 0 \leq r<4,0<t, y \in C_{t}$.

This is proved by straightforward estimates (and the Ito formula) using each time $-\frac{1}{8} \int_{0}^{t} x_{s}^{6} d s$ to keep the contributions of each of the other terms in check.
4.8. A robust version for $\oint_{t}=E\left[\phi\left(x_{t}\right) \mid Y_{t}\right]$

Now let $B$ denote the space of bounded Borel functions $\phi: \mathbb{R} \rightarrow \mathbb{R}$ endowed with the sup norm. Define for $\phi \in B$

$$
\begin{equation*}
\left\langle N_{t}(y), \phi\right\rangle=E\left[\phi\left(x_{t}\right) U(0,0, t, y)\right] \tag{4.9}
\end{equation*}
$$

This is well defined because $U(0,0, t, y) \in L^{1}(\Omega, A, P)$ and $\phi\left(x_{t}\right)$ is bounded. One now shows using 4.7 that the finite positive measure $N_{t}(y)$ has a density $n_{t}(y)(x) \in L^{1}(\mathbb{R})$. Now define for all $\phi \in E$

$$
\begin{equation*}
\left\langle N_{t}(y), \phi\right\rangle=\int_{-\infty}^{+\infty} \phi(x) n_{t}(y)(x) d x \tag{4.10}
\end{equation*}
$$

(which agrees (4.9) for $\phi \in B$ ). The functional $N_{t}(y): E \rightarrow \mathbb{R}$ satisfies

$$
\begin{equation*}
\left\|<N_{t}(y), \phi>\right\| \leq c\|\phi\|_{(\beta, r)} \quad \text { for } \phi \in \varepsilon(\beta, r) \tag{4.11}
\end{equation*}
$$

with $C$ independent of $y$, $t$ for bounded $\|y\|_{t}$ and $t$ in a compact subset of ( $0, \infty$ ).
4.12. Proposition. $N_{t}(y)$ depends continuously on $y$, for each fixed $t$, $\phi \in E$.

$$
\text { Now define } \Delta_{t}^{\phi}: C_{t} \rightarrow \mathbb{R} \text { by }
$$

(4.13) $\quad \Delta_{t}^{\phi}(y)=\left\langle N_{t}(y), \phi\right\rangle$.

Then $\Delta_{t}^{\phi}(y)$ is continuous and one shows that $\Delta_{t}^{1}(y)>0$ for all $t>0$ so that we can define
(4.14) $\quad \delta_{t}^{\phi}(y)=\frac{\Delta_{t}^{\phi}(y)}{\Delta_{t}^{1}(y)}$.
4.15. Proposition. If $\phi \in E, \omega \mapsto \delta_{t}^{\phi}\left(y_{t}(\omega)\right)$, where $y_{t}(\omega)$ is the $\omega$-path of the process $y$ of (2.1) is a version of $\phi_{t}=\mathrm{E}\left[\phi\left(\mathrm{x}_{\mathrm{t}}\right) \mid \mathrm{y}_{\mathrm{t}}\right]$ for $\mathrm{t}>0$. (So $\delta_{t}^{\phi}$ is a robust pathwise version of $\phi_{t}$.)

This is proved via the Kallianpur-Stiebel formula which says in our case that if

$$
k_{t}(\phi)=\int_{\Omega} \phi\left(x_{t}\left(\omega^{\prime}\right)\right) \exp \left(\int_{0}^{t} x_{s}^{3}\left(\omega^{\prime}\right) d y_{s}(\omega)-\frac{1}{2} \int_{0}^{t} x_{s}^{6}\left(\omega^{\prime}\right) d s\right) d P\left(\omega^{\prime}\right)
$$

then

$$
\phi_{t}=\frac{k_{t}(\phi)}{k_{t}(1)}
$$

Indeed rewrite the term $\int_{0}^{t} x_{s}^{3}\left(\omega^{\prime}\right) d y_{s}(\omega)$ by means of partial integration as

$$
\begin{aligned}
\int_{0}^{t} x_{s}^{3}\left(\omega^{\prime}\right) d y_{s}(\omega) & =x_{t}^{3}\left(\omega^{\prime}\right) y_{t}(\omega)-\int_{0}^{t} y_{s}(\omega) d\left(x_{s}^{3}\left(\omega^{\prime}\right)\right)= \\
& =x_{t}^{3}\left(\omega^{\prime}\right) y_{t}(\omega)-3 \int_{0}^{t} y_{s}(\omega) x_{s}^{2}\left(\omega^{\prime}\right) d w_{s}\left(\omega^{\prime}\right)- \\
& -3 \int_{0}^{t} y_{s}(\omega) x_{s}(\omega) d s
\end{aligned}
$$

and it readily follows that $\int U\left(0,0, t, y_{t}(\omega)\right) \phi\left(x_{t}\left(\omega^{\prime}\right)\right) d P\left(\omega^{\prime}\right)=\Delta_{t}^{\phi}(y)$ is a version of $k_{t}(\phi)$.
4.16. Smoothness properties of $n_{t}(y)(x)$

Let $F$ be the space of all $C^{\infty}$-functions $\phi: \mathbb{R} \rightarrow \mathbb{R}$ for which
$\exp \left(\beta|x|^{r}\right)\left|\phi^{(k)}(x)\right|$ is bounded for all $\beta \geq 0,0 \leq r<4, k \in \mathbb{N} \cup\{0\}$ and give $F$ the topology defined by the family of norms
(4.17) $\quad\|\phi\|_{F(\beta, r, k)}=\sup \left\{\exp \left(\beta|x|^{r}\right)\left|\phi^{(k)}(x)\right|\right\}$.
4.18. Lemma. If $x_{0} \in H$ then $n_{t}(y) \in F$ for $a Z Z t>0, y \in C_{t}$ and $n_{t}(y)(x)>0$ for all $\mathrm{x} \in \mathbb{R}, \mathrm{t}>0$ for y differentiable.

This is approached by considering the derivatives of $N_{t}(y)$ defined by $\left\langle N_{t}(y)^{\prime}, \phi\right\rangle=\left\langle N_{t}(y), \phi^{\prime}\right\rangle$ for smooth $\phi: \mathbb{R} \rightarrow \mathbb{R}$.

### 4.19. Robustness for the filter

Now consider a stochastic differential equation with output map and initial condition driven by the observation process $y_{t}$

$$
\begin{equation*}
d z=\alpha(z) d t+\beta(z) d y_{t}, \quad z(0)=z^{\text {in }}, \quad z \mapsto \gamma(z), \quad z \in M \tag{4.20}
\end{equation*}
$$

as we would have for a filter for $\phi$ cf. 2.4 above. Equation (4.20) is to be interpreted in the Stratonovic sense. Let $T$ be the stopping time for a maximal solution. Then, as was shown in [Sussman, 1978] these equations admit robust solutions in the following sense.

Consider the equation for $y \in C_{t}$
(4.21) $d z=\alpha(z) d t+\beta(z) d y, \quad z(0)=z^{\text {in }}$.

A curve $z: \tau \rightarrow z(\tau), 0 \leq \tau \leq t$ is said to be a solution of (4.21) if there exists a neighborhood $U$ of $y$ in $C_{t}$ with the property that there is a continuous map $U \rightarrow C([0, t], M) \tilde{y} \rightarrow z(\tilde{y})$ to the space of continuous curves in $M$ such that $z(\tilde{y})$ is a solution of (4.21) in the usual sense for all $\tilde{y} \in U \cap C_{t}^{l}$ (so that the equation can be written as a usual differential
equation) and $z(y)=z$.
With this notion of solutions the robustness result is:
4.22. THEOREM [Sussman, 1978]. (i) Given any continuous $y:[0, \infty) \rightarrow \mathbb{R}$, $\mathrm{y}(0)=0$, there exists a time $\mathrm{T}(\mathrm{y})>0$ such that there is a unique solution $\tau \rightarrow z(y)(\tau)$ of (4.21). If $T(y)<\infty$ then $\{z(y)(t): 0 \leq t<T(y)\}$ is not relatively compact on $M$.
(ii) If y is a Wiener process with continuous sample paths defined on $(\Omega, A, P)$ and if $y^{\omega}(t)=y_{t}(\omega)$, then $\omega \rightarrow T\left(y^{\omega}\right)$ is a version of the stopping time up to which the Stratonovic solution of (4.20) is defined and $\omega \rightarrow z\left(y^{\omega}\right)(t), 0 \leq t<T\left(y^{\omega}\right)$ is a version of the solution $z_{t}$ for each $t>0$.

In our setting $y_{t}$ is not a Wiener process, but the same techniquess apply, and the same results hold.

In other words up to a stopping time, solutions of (4.20) exist pathwise, they are continuous as a function of the path and hence can be calculated as limits of solutions to the corresponding nonstochastic differential equations (4.21) for (piecewise) differential continuous y.
23. Everywhere equality of the robust filter output and the robust DMZ output and consequences

Now let (4.20) be a smooth filter for $\phi$ in the sense of section 2.4 $\geq$ and let $\phi \in E$. Choose the robust version of $\phi_{t}$, i.e. the map $\varsigma_{t}^{\phi}\left(y^{\omega}\right)$ and choose the robust solution $\omega \rightarrow z\left(y^{\omega}\right)(t)$ of (4.20). The fact (4.20) is a filter for $\phi_{t}$ says by definition that $\delta_{t}^{\phi}\left(y^{\omega}\right)=\gamma\left(z\left(y^{\omega}\right)(t)\right)$ almost all $\omega$ such that $T\left(y^{\omega}\right)>t$. The robustness of the two versions ceadily implies that
) $\quad \delta_{t}^{\phi}(y)=\gamma(z(y)(t))$
s everywhere whenever $t>0, y \in C_{t}, T(y)>t$.
4.25. Smoothness properties of the family of densities $n_{t}(y)$

When $y$ is piecewise $C^{1}$ and the initial probability density $\nu^{i n}$ is in $F$ the study of the measure $N_{t}(y)$ is much easier. (By modifying the data $(\Omega, A, P)$, etc. it can actually be arranged that $\nu^{i n}$ is in $F$ essentially by replacing $\nu^{\text {in }}$ with the density at a slightly later time $\tau<t$ ). Now integration by parts gives

$$
\begin{equation*}
u(0,0, t, y)=\exp \left(\int_{0}^{t} x_{s}^{3} \dot{y}(s) d s-\frac{1}{2} \int_{0}^{t} x_{s}^{6} d s\right) \tag{4.26}
\end{equation*}
$$

Now let $\phi: \mathbb{R} \rightarrow \mathbb{R}$ be of class $C^{2}$, then the differential of $\phi\left(x_{t}\right) U(0,0, t, y)$ can easily be computed to be

$$
\begin{equation*}
\left[\phi^{\prime}\left(x_{t}\right) d w_{t}+\frac{1}{2} \phi^{\prime \prime}\left(x_{t}\right) d t+\left(x_{t}^{3} \dot{y}(t)-\frac{1}{2} x_{t}^{6}\right) \phi\left(x_{t}\right)\right] u(0,0, t, y) \tag{4.27}
\end{equation*}
$$

so that (if, say, $\phi$ has compact support), ${ }_{T}$ denoting the restriction of $y$ to $[0, \tau]$

$$
\begin{align*}
& E\left[\phi\left(x_{t}\right) U(0,0, t, y)-\phi\left(x_{\tau}\right) U\left(0,0, \tau, y_{\tau}\right)\right]=  \tag{4.28}\\
& \left.\int_{\tau}^{t}\left(\frac{1}{2} \frac{d^{2}}{d x^{2}}+x^{3} \dot{y}(t)-\frac{x^{6}}{2}\right) \phi(x)\right)\left.\right|_{x=x_{s}} U\left(0,0, s, y_{s}\right) d s
\end{align*}
$$

and this in turn says that the densities $n_{t}(y)$ of $N_{t}(y)$, i.e. the functions $(t, x) \rightarrow n_{t}(y)(x)$ satisfy the partial differential equation

$$
\begin{equation*}
\frac{\partial \rho}{\partial t}=\frac{1}{2} \frac{\partial^{2} \rho}{\partial x^{2}}+\left(x^{3} u(t)-\frac{x^{6}}{2}\right) \rho, \quad \rho(0, x)=n_{0}(x) \tag{4.29}
\end{equation*}
$$

where $n_{0}$, the initial density is in $F$ and $u=\dot{y}$. One has
4.30. LEMMA. Let $u$ be piecewise continuous an $[0, T]$, and for each $n_{0} \in F$ Let $\rho_{n_{0}, t}$ be the function $x \rightarrow \rho(t, x)$ where $\rho$ solves (4.29) then $\left(n_{0}, t\right) \rightarrow \rho_{n_{0}, t}$ : $F \times[0, T] \rightarrow F$ is continuous.

Now let

$$
\begin{equation*}
L_{0}=\frac{1}{2} \frac{d^{2}}{d x^{2}}-\frac{x^{6}}{2}, \quad L_{1}=x^{3} \tag{4.31}
\end{equation*}
$$

considered as (differential) operators $F \rightarrow F$. For each constant $\bar{u}$ let $L(\bar{u})=L_{0}+\bar{u}_{1}$ and let $\exp (t L(\bar{u})) \psi$ for $\psi \in F$ denote the solution of (4.29) with $u(\cdot) \equiv \bar{u}, n_{0}=\psi$.

Let $K \subset \mathbb{R}^{n}$ be a convex subset with nonempty interior. A family $\{\psi(v): v \in K\}$ of elements of $F$ is said to depend smoothly on $v$ if $(x, v) \rightarrow \psi(v)(x)$ is a $C^{\infty}$ function on $K \times \mathbb{R}$ and for each $\underline{m}=\left(m_{1}, \ldots, m_{n}\right)$, $\mathrm{v} \mapsto\left(\frac{\partial^{\underline{m}}}{\partial \mathrm{v}^{\underline{m}}} \psi\right)(\mathrm{v}), \quad \mathrm{v} \in \mathrm{K}$ takes values in $F$ and is a continuous map $K \rightarrow F$. One then has
4.32. LEMMA. $\{\exp (\mathrm{tL}(\overline{\mathrm{u}})) \psi(\mathrm{v})\}$ depends smoothly on $(\mathrm{v}, \mathrm{t})$ if $\psi(\mathrm{v}), \mathrm{v} \in \mathrm{K}$ is a smooth family.
4.33. Corollary. Let $\overline{\mathrm{u}}_{1}, \ldots, \overline{\mathrm{u}}_{\mathrm{m}} \in \mathbb{R}$. Then if $\psi \in \mathcal{F}$ the fomily $\left\{\exp \left(\mathrm{t}_{1} \mathrm{~L}\left(\bar{u}_{1}\right)\right) \ldots \exp \left(\mathrm{t}_{\mathrm{m}} \mathrm{L}\left(\bar{u}_{\mathrm{m}}\right)\right)^{m} \psi:\left(\mathrm{t}_{1}, \ldots, \mathrm{t}_{\mathrm{m}}\right) \in[0, \infty]^{m^{m}}\right\}$ depends continuous $l_{y}$ on $t_{1}, \ldots, t_{m}$; moreover for each $\underline{\mu}=\left(\underline{\mu}_{1}, \ldots, \underline{\mu}_{m}\right)$ we have

$$
\begin{align*}
& \frac{\partial^{\underline{\underline{u}}}}{\partial t^{\underline{\mu}}}\left(\exp \left(t_{1} L\left(\bar{u}_{1}\right)\right) \ldots \exp \left(t_{m} L\left(\bar{u}_{m}\right)\right) \psi\right)=  \tag{4.34}\\
& L\left(\bar{u}_{1}\right)^{\mu_{1}} \exp \left(t_{1} L\left(\bar{u}_{1}\right)\right) L\left(\bar{u}_{2}\right)^{\mu_{2}} \exp \left(t_{2} L\left(\bar{u}_{2}\right)\right) \ldots L\left(\bar{u}_{m}\right)^{\mu}{ }_{m} \exp \left(t_{m} L\left(\bar{u}_{m}\right)\right) \psi .
\end{align*}
$$

SYSTEM THEORETIC PART III: REALIZATION THEORY

## 1. Some differential topology on $F$

Let $U$ be an open subset of the space of smooth functions $F$. A map $U \rightarrow \mathbb{R}$ is said to be of class $C^{\infty}$ if the function $\nu \rightarrow \Phi(\psi(\nu))$ is $C^{\infty}$ in the ual sense for every family $\{\psi(v): v \in K\}$ depending smoothly on $v$ in the :nse described in section 4 above. This class of functions is denoted $C^{\infty}(U)$. $\lambda$ is a continuous linear functioned on $F$ then $\lambda$ (restricted to any $U$ ) is : class $C^{\infty}$. Note that $C^{\infty}(U)$ is closed under pointwise multiplication and ivision by functions in $C^{\infty}(U)$ which are everywhere nonzero.

Let $L$ be a continuous linear operator on $F$, then $L$ defines a "linear vectorfield"
$\tilde{L}: C^{\infty}(F) \rightarrow C^{\infty}(F) \quad$ (and $C^{\infty}(U) \rightarrow C^{\infty}(U)$ for each $U$ ) defined by

$$
\begin{equation*}
(\tilde{L} \Phi)(\psi)=\left.\frac{d}{d t}\right|_{t=0} \Phi(\psi+t L \psi) \tag{5.2}
\end{equation*}
$$

This is completely analogous to the map which assigns to an $n \times n$ matrix $A=\left(a_{i j}\right)$ the "linear vector field" $\Sigma a_{i j} x_{i} \frac{\partial}{\partial x_{j}}$. It is totally routine to
check that check that

$$
\begin{equation*}
\left[\mathrm{L}_{1}, \mathrm{~L}_{2}\right]^{\sim}=-\left[\tilde{\mathrm{L}}_{1}, \tilde{\mathrm{~L}}_{2}\right] \tag{5.3}
\end{equation*}
$$

5.4. LEMMA. Let $\{\psi(t): 0 \leq t<\varepsilon\} \subset U$ depend smoothly on $t$ and let $\dot{\psi}$ be the $t$-derivative of $\psi$. Then for all $\Phi \in C^{\infty}(U)$

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{dt}}\right|_{\mathrm{t}=0} \Phi(\psi(\mathrm{t}))=\frac{\mathrm{d}}{\mathrm{dt}} \Phi(\psi(0)+\mathrm{t} \dot{\psi}(0)) \tag{5.5}
\end{equation*}
$$

In particular if $L$ is a continuous linear operator on $F$ such that $\mathrm{L} \psi(0)=\dot{\psi}(0)$, then

$$
\begin{equation*}
(\tilde{\mathrm{L}} \Phi)(\psi(0))=\left.\frac{\mathrm{d}}{\mathrm{dt}}\right|_{\mathrm{t}=0} \Phi(\psi(\mathrm{t})) . \tag{5.6}
\end{equation*}
$$

Now let $U \subset F$ be the set of all $\psi \in F$ such that $\int_{-\infty}^{+\infty} \psi(x) d x>0$ and let $\Phi: U \rightarrow \mathbb{R}$ be given by the kind of formula occuring in our conditional expectation expressions
(5.7) $\quad \Phi(\psi)=\frac{\int \phi(x) \psi(x) d x}{\int \psi(x) d x}$.

For the smooth families $e^{\operatorname{tL}(\overline{\mathrm{u}})} \psi$, where $L(\overline{\mathrm{u}})$ is as in 4.32 above, one finds (5.8) $\quad(\tilde{\mathrm{L}}(\overline{\mathrm{u}}) \Phi)(\psi)=\frac{\mathrm{d}}{\mathrm{dt}} \Phi\left(\mathrm{e}^{\mathrm{tL}(\overline{\mathrm{u}})} \psi\right)$ and repeating this

$$
\begin{align*}
& \tilde{L}\left(\bar{u}_{m}\right) \ldots \tilde{L}\left(\bar{u}_{1}\right) \Phi(\psi)=  \tag{5.9}\\
& \left.\frac{\partial^{m}}{\partial t_{1} \partial t_{2} \ldots \partial t_{m}}\right|_{1}=\ldots=t_{m}=0^{\Phi\left(e^{t_{1} L\left(\bar{u}_{1}\right)} \ldots e^{t_{m} L\left(\bar{u}_{m}\right)} \psi\right) .}
\end{align*}
$$

### 5.10. The Lie-algebraic inplications of the existence of a smooth filter

Now let us repeat these remarks for the more familiar case of vectorfields $\alpha, \beta$ on a smooth finite dimensional manifold $M$ with for each $\bar{u} \in \mathbb{R}, A(\bar{u})$ the vectorfield $\alpha+\bar{u} \beta$. Let $\left.\pi\left(\left(\bar{u}_{1}, t_{1}\right),\left(\bar{u}_{2}, t_{2}\right), \ldots,\left(\bar{u}_{m}, t_{m}\right)\right) ; z\right)$ be the result of letting $z$ evolve on $M$ along $A\left(\bar{u}_{m}\right)$ during time $t_{m}$, then along $A\left(\bar{u}_{m-1}\right)$ during time $t_{m-1}, \ldots$ Let $\gamma: M \rightarrow \mathbb{R}$ be a smooth function. We have of course

$$
\begin{equation*}
(A(\bar{u}) \gamma)(z)=\left.\frac{d}{d t}\right|_{t=0} \gamma(\pi((\bar{u}, t) ; z)) \tag{5.11}
\end{equation*}
$$

and

$$
\begin{align*}
& \left(A\left(\bar{u}_{m}\right) \ldots A\left(\bar{u}_{1}\right) \gamma\right)(z)=  \tag{5.12}\\
& \left.\frac{\partial^{m}}{\partial t_{1} \ldots \partial t_{m}} \right\rvert\, t_{1}=\ldots=t_{m}=0 \quad \gamma\left(\pi\left(\left(\bar{u}_{1}, t_{1}\right), \ldots,\left(\bar{u}_{m}, t_{m}\right) ; z\right) .\right.
\end{align*}
$$

Let $R \subset M$ be the set of all points in $M$ which can be reached from $z^{\text {in }}$ $v$ means of these bang-bang-bang controls in time $<\bar{T}$ i.e. $R$ is the set of $L \pi\left(\left(\bar{u}_{1}, t_{1}\right), \ldots,\left(\bar{u}_{m}, t_{m}\right) ; z^{i n}\right)$ with $\sum t_{i}<\bar{T}$. Let $z \in R$ and choose a bangag control which steers $z^{\text {in }}$ to $z$ in time $\tau<\bar{T}$; let $\psi_{z} \in F$ be the solution the "control version" of the $O M Z$ equation (4.29), with initial condition the density of $\nu^{\text {in }}$. Then $\psi_{z} \in U$, because $\Delta_{t}^{1}(y)>0$ (cf. just below position 4.12). Now let $\bar{u}_{1}, \ldots, \bar{u}_{m}, t_{1}, \ldots, t_{m}$ satisfy $\left|\bar{u}_{i}\right|=1$, $\left|, \ldots, t_{m}\right|<\bar{T}-\tau, t_{i} \geq 0$ and assume that $(\alpha, \beta, \gamma)$ on $M$ define a smooth lter for a given $\phi \in E$ in the sense of subsection 2.4. Let $\Phi$ be the nresponding functional (5.7). Then by (4.24) we have
i.13) $\left.\Phi\left(e^{t L_{1}\left(\bar{u}_{1}\right)} \ldots e^{t m_{m}^{L\left(\bar{u}_{m}\right)}} \psi_{z}\right)=\gamma\left(\pi\left(\left(\bar{u}_{1}, t_{1}\right)\right), \ldots,\left(\bar{u}_{m}, t_{m}\right)\right) ; z\right)$
and this was really the whole reason for establishing formula (4.24), that is the reason why we needed to prove the existence of a robust pathwise version $\hat{\phi}_{t}$ ).

Now let $\underset{\underline{A}}{ }$ denote the free associative algebra on two generators $a_{-}, a_{+}$Let $\underset{=}{A}$ be the associative algebra (under composition) of linear maps
$C^{\infty}(U) \rightarrow C^{\infty}(U)$ generated by $\tilde{L}(-1), \tilde{L}(1)$ and let ${\underset{=}{A}}_{2}$ be the associative algebra (again under composition) of differential operators on M. Homomor-
 $v_{1}\left(a_{-}\right)=\tilde{L}(-1), v_{1}\left(a_{+}\right)=\tilde{L}(1), v_{2}\left(a_{-}\right)=A(-1), v_{2}\left(a_{+}\right)=A(1)$.

Let $\xlongequal[\equiv]{ }$ denote the free Lie-algebra on the generators $a_{-}, a_{+}$(viewed as a subalgebra of A) and let $\underline{\underline{L}}_{1}$, $\underline{\underline{L}}_{2}$ denote the Lie algebras generated respectively by $\tilde{L}(-1), \tilde{L}(1)$ and $A(-1), A(1)$ (as subalgebras of ${\underset{A}{A}}$ and ${\underset{N}{2}}^{2}$ ). Then of course we have induced homomorphisms $v_{i}: L \rightarrow L_{i}$
(5.14)


Let I denote the set of those vectorfields $V \in \underline{\underline{I}}_{2}$ such that

$$
\begin{equation*}
\left(\left[\mathrm{v}_{1},\left[\mathrm{v}_{2},\left[\ldots\left[\mathrm{v}_{\mathrm{m}}, \mathrm{v}\right] \ldots\right]\right] \gamma\right)(\mathrm{z})=0 \text { for all } \mathrm{v}_{1}, \ldots, \mathrm{v}_{\mathrm{m}} \in{\stackrel{L}{L_{2}}}_{2}\right. \tag{5.15}
\end{equation*}
$$

5.16. LEMMA. I is an ideal and if $a \in \underline{I}$ is such that $\nu_{1}(a)=0$ then $\nu_{2}(a) \in I$.

It follows that there is a homomorphism of Lie algebras $L_{1} \rightarrow L_{2} / I$ making diagram (5.14) commutative.
The lemma is proved by combining (5.12) and (5.9) with (5.13) and this is why we needed to establish smoothness properties of families like $\exp (\mathrm{tL}(\overline{\mathrm{u}}) \psi)$.

### 5.17. Foliations and such

The last step in this section is now to show that $\underline{L}_{2} / I$, or more precisely a suitable quotient, is (isomorphic to) a subalgebra of a Lie algebra of vectorfields in a smooth finite dimensional manifold (a subquotient manifold of $M$ ) for suitable $z$. Let $S$ be a set of smooth vectorfields on $M$. For each $z \in M$ consider $S(z)=\{V(z): V \in S\}$. For $\underline{L}_{2}$ and $I, L_{2}(z)$ and $I(z)$ are vectorspaces. Let $k$ be the maximum of the $\operatorname{dim} \underline{L}_{2}(z)$ for $z \in R$ and $k_{0}$ the maximum of the $\operatorname{dim} I(z)$, for $z \in R$ and $\operatorname{dim} \underline{L}_{2}(z)=k$. Choose a $\bar{z}$ in the relative interior of $R$ such that $\operatorname{dim} \underline{L}_{2}(\bar{z})=k, \operatorname{dim} I(\bar{z})=k_{0}$, and choose a neighborhood $N$ of $\bar{z}$ (in M) such that $\operatorname{dim}{\underset{\underline{L}}{2}}^{\underline{L}_{2}}(z) \geq k, \operatorname{dim} I(z) \geq k_{0}$ for $z \in N$. (This can be done (obviously)). Suppose $\bar{z}$ can be reached from $z^{\text {in }}$ in time $\tau<\bar{T}$. Now let $M_{0}$ be a connected submanifold of $N$ of which all points are reachable from $\bar{z}$ in time $\langle\bar{T}-\tau-\delta$ for some $\delta>0$ and which is maximal in dimension in the set of all such manifolds. Then $\alpha, \beta$ are tangent to $M_{0}$ and $\bar{z}$ (as is easily checked). Also (becaume $M_{0} \subset N$ and $M_{0} \subset R$ ) we have $\operatorname{dim}\left(\underline{\underline{L}}_{2}(z)\right)=k$ for all $z \in M_{0}$ so that by Frobenius theorem there exists a submanifold $M_{1}$ of $M_{0}$ whose tangent space at each point $z \in M_{1}$ is precisely $\underline{I}_{2}(z)$. One then also has that $\operatorname{dim} I(z)=k_{0}$ for all $z \in M_{1}$ so that $I$ has integral manifolds $M_{2}$ locally near $\bar{z}, M_{1}$ is then foliated by the integral manifolds of $I$ so that $M_{1}$ locally near $\bar{z}$ looks like $M_{1} \simeq M_{2} \times M_{3}$. The Lie algebra of vectorfields of $M_{3}$ is then isomorphic to the quotient $\underline{L}_{2} \mid M_{1} / I_{\mid M}$. So That by restriction to $\mathrm{M}_{1}$ the dotted arrow is diagram 5.14 gives a homomorphism of Lie algebras

$$
\begin{equation*}
\stackrel{L}{=}_{1} \rightarrow \mathrm{~V}\left(\mathrm{M}_{3}\right) \tag{5.18}
\end{equation*}
$$

5.19. Proposition. Assume that the homomorphism of Lie algebras (5.18) is zero and assume moreover that $\underline{L}_{1}$ contains all the operators $\mathrm{L}_{\mathrm{k}}=\frac{\mathrm{d}}{\mathrm{dx}} \mathrm{x}^{\mathrm{k}}$, $\mathrm{k}=0,1, \ldots$. Then $\phi$ is a constant almost everywhere.

This is seen as follows. This homomorphism is zero iff $k_{0}=k$ so that for $z \in R, V \gamma(z)=0$ for all $V \in \underline{L}_{2}$, which gives $(\tilde{L} \Phi)\left(\psi_{2}\right)=0$ for all $L \in \underline{L}_{2}$. Now calculate $(\tilde{L} \Phi)\left(\psi_{2}\right)$ using formula (5.7) for $\Phi\left(\psi_{2}\right)$ to find that

$$
\begin{equation*}
\left.\left.<\phi, L \psi_{z}\right\rangle<1, \psi_{z}\right\rangle=\left\langle\phi, \psi_{z}\right\rangle\left\langle 1, \mathrm{I} \psi_{z}\right\rangle \text { for all } \mathrm{L} \in \underline{\mathrm{~L}}_{2} . \tag{5.20}
\end{equation*}
$$

As $<1, L_{k} \psi_{z}>=0$ this gives

$$
\begin{equation*}
\int \phi(x)\left[x^{k} \psi_{z}^{\prime}(x)+k x^{k-1} \psi_{z}(x)\right] d x=0, \quad k=0,1,2, \ldots \tag{5.21}
\end{equation*}
$$

From this, using that $\phi(x) \psi_{z}^{\prime}(x)$ and $\phi(x) \psi_{2}(x)$ are bounded by $e^{\beta|x|^{r}}$ for some $\beta>0, r<4$, one sees by considering the Fourier transforms of $\phi(x) \psi_{z}^{\prime}(x)$ and $\phi(x) \psi_{z}(x)$ that $\phi^{\prime} \psi_{z}=0$ and as $\psi_{2}$ never vanishes that $\phi$ is constant.
6. ALGEBRAIC PART
6.1. The Weyl Lie algebras $W_{n}$

The Weyl Lie algebra $W_{n}$ is the algebra of all differential operators (any order) in $\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n}}$ with polynomial coefficients. The Lie brackett operation is of course the commutator $\left[D_{1}, D_{2}\right]=D_{1} D_{2}-D_{2} D_{1}$. A basis for $W_{1}$ (as a vector space over $\mathbb{R}$ ) constists of the operators
(6.2) $\quad x^{i} \frac{\partial^{j}}{\partial x^{j}}, i, j=0,1,2, \ldots$
(where of course $x^{i} \frac{\partial^{0}}{\partial x^{0}}=x^{i}, x^{0} \frac{\partial^{j}}{\partial x^{j}}=\frac{\partial^{j}}{\partial x^{j}}, x^{0} \frac{\partial^{0}}{\partial x^{0}}=1$ ). One has for example

$$
\left[\frac{\partial^{2}}{\partial x^{2}}, x^{2}\right]=4 x \frac{\partial}{\partial x}+2
$$

as is easily verified by calculating

$$
\left[\frac{\partial^{2}}{\partial x^{2}}, x^{2}\right] f(x)=\frac{\partial^{2}}{\partial x^{2}}\left(x^{2} f(x)\right)-x^{2} \frac{\partial^{2}}{\partial x^{2}}(f(x))
$$

for an arbitrary test function (polynomial) $f(x)$.
Some easy facts (theorems) concerning the Weyl Lie algebras $W_{n}$ are (cf. [Hazewinkel-Marcus, 1981] for proofs):
6.3. Proposition. The Lie algebra $W_{n}$ is generated (as a Lie algebra) by the elements $x_{i}, \partial^{2} / \partial x_{i}^{2}, x_{i}^{2}\left(\partial / \partial x_{i}\right), i=1, \ldots, n ; x_{i} x_{i-1}, i=2, \ldots, n$. In particular $W_{1}$ is generated by $x, \partial^{\frac{2}{2}} / \partial x^{2}, x^{2}(\partial / \partial x)$.
6.4. Proposition. The only nontrivial ideal of $W_{n}$ is the one-dimensional ideal $\mathbb{R} 1$ of scalar multiples of the identity operator.

If $M$ is a $C^{\infty}$ differentable manifold let $V(M)$ denote the Lie-algebra of all $C^{\infty}$ vectorfields on $M$ (i.e. the Lie algebra of all derivations on the ring of smooth functions on $M$ ). If $M=\mathbb{R}^{n}, V\left(\mathbb{R}^{n}\right)$ is the Lie algebra of all differential operators of the form

$$
\sum_{i=1}^{n} g_{i}\left(x_{1}, \ldots, x_{n}\right) \frac{\partial}{\partial x_{i}}
$$

with $g_{i}\left(x_{1}, \ldots, x_{n}\right)$ a smooth function on $\mathbb{R}^{n}$.
A deep fact concerning the Weyl Lie algebras $W_{n}$ is now
6.5. Theorem. Let $M$ be a finite dimensional smooth manifold. Then there are no nonzero homomorphisms of Lie algebras $W_{n} \rightarrow V(M)$ or $W_{n} / \mathbb{R} 1 \rightarrow V(M)$ for $\mathrm{n} \geq 1$.

The original proof of this theorem ([Hazewinkel-Marcus, 1981] is long and computational. Fortunately there now exists a much better proof (about two pages) of the main and most difficult part [Stafford, 1982], essentially based on the observation that the associative algebra $W_{1}$ cannot have left ideals of finite codimension. For some more remarks about the proof cf. 6.8 below.
6.6. The Lie algebra of the cubic sensor

According to section 2 above the estimation Lie algebra $L(\Sigma)$ of the cubic sensor is generated by the two operators

$$
L_{0}=\frac{1}{2} \frac{d^{2}}{d x^{2}}-\frac{1}{2} x^{6}, \quad L_{1}=x^{3}
$$

Calculating $\left[L_{0}, L_{1}\right]$ gives $C=3 x^{2} \frac{d}{d x}+3 x$. Let $\operatorname{ad}_{C}(-)=[C,-]$. Then
$\left(\operatorname{ad}_{C}\right)^{3} B=C^{s t} t_{x}^{6}$ which combined with A gives as that $\left(d^{2} / d x^{2}\right) \in L(\Sigma) \cdot$ To show that also $x^{2} \frac{d}{d x} \in L(\Sigma)$ requires the calculation of some more bracketts (about 15 of them). For the details cf. [Hazewinkel-Marcus, 1981]. Then $\mathrm{x}, \mathrm{x}^{2} \frac{\mathrm{~d}}{\mathrm{dx}}, \frac{\mathrm{d}^{2}}{\mathrm{dx}} \in \mathrm{L}(\bar{z})$ which by proposition 6.3 implies:
6.7. Theorem. The estimation Lie algebra $L(\Sigma)$ of the cubic sensor is equal to the Weyl Lie algebra $\mathrm{W}_{1}$.

In a similar manner one can e.g. show that the estimation lie algebra of the system $d x_{t}=d w_{t}, d y_{t}=\left(x_{t}+\varepsilon x_{t}^{3}\right) d t+d v_{t}$ is equal to $W_{1}$ for all $\varepsilon \neq 0$. It seems highly likely that this is a generic phenomenon i.e. that the estimation Lie algebra of a system of the form $d x_{t}=f\left(x_{t}\right) d t+G\left(x_{t}\right) d t$, $d y_{t}=h\left(x_{t}\right) d t+d v_{t}$ with $x \in \mathbb{R}^{n}$ and $f, G$ and $h$ polynomial is equal to $W_{n}$ for almost all (in the Zariski topology sense) polynomials f, G, h.

### 6.8. Outline of the proof of the nonembedding theorem 6.5

Let $\hat{\mathrm{V}}_{\mathrm{n}}$ be the Lie algebra of all expressions
(6.9)

$$
\sum_{i=1}^{n} f_{i}\left(x_{1}, \ldots, x_{n}\right) \frac{\partial}{\partial x_{i}}
$$

where $f_{1}(x), \ldots, f_{n}(x)$ are formal power series in $x_{1}, \ldots, x_{n}$. (No convergenc properties are required). Suppose that
(6.10) $\alpha: W_{n} \rightarrow V(M)$
is a nonzero homomorphism of Lie algebras into some $V(M)$ with $M$ finite dimensional. Then there is $a D \in W_{n}$ and an $m \in M$ such that the tangent vector $\alpha(D)(m) \neq 0$. Now take formal Taylor series of the $\alpha(D)$ around $m$ (with respect to local coordinates at $m$ ) to find a nonzero homomorphism of Lie algebra
(6.11) $\quad \hat{\alpha}: W_{n} \rightarrow \hat{V}_{m}$
where $m=\operatorname{dim}(M)$.

Observe that $W_{1}$ is a sub-algebra of $W_{n}$ (consisting of all differential operators not involving $x_{i}, i \geq 2$, and $\partial / \partial x_{i}, i \geq 2$ ) so that $i t$ suffices to prove theorem 6.5 for the case $n=1$.

Because the only nontrivial ideal of $W_{1}$ is $\mathbb{R} 1$ (cf. proposition 6.4) the existence of a nonzero $\hat{\alpha}: W_{1} \rightarrow \hat{V}_{m}$ implies that $W_{1}$ or $W_{1} / \mathbb{R} 1$ can be embedded in $\hat{V}_{\mathrm{m}}$.

The Lie-algebra $\hat{\mathrm{V}}_{\mathrm{m}}$ carries a filtration $\hat{\mathrm{V}}_{\mathrm{m}}=\mathrm{L}_{-1} \supset \mathrm{~L}_{0} \supset \mathrm{~L}_{1} \supset \ldots$ where the $L_{i}$ are sub-Lie-algebras. This filtration has the following properties

$$
\begin{equation*}
\left[L_{i}, L_{j}\right] \subset\left[L_{i+j}\right] \tag{6.12}
\end{equation*}
$$

$$
\begin{equation*}
\bigcap_{i=-1}^{\infty} L_{i}=\{0\} \tag{6.13}
\end{equation*}
$$

$$
\begin{equation*}
\operatorname{dim}\left(L_{-1} / L_{i}\right)<\infty, \quad i=-1,0,1, \ldots \tag{6.14}
\end{equation*}
$$

where "dim" means dimension of real vectorspaces.
Indeed let

$$
\begin{equation*}
f_{i}\left(x_{1}, \ldots, x_{n}\right)=\sum_{v} a_{i, v} x^{v} \tag{6.15}
\end{equation*}
$$

$\nu=\left(\nu_{1}, \ldots, \nu_{m}\right), \nu_{i} \in \mathbb{N} \cup\{0\}$ a multi index, be the explicit power series for $f_{i}(x)$. Then $L_{j} \subset \hat{V}_{m}$ consists of all formal vectorfields (6.15) for which

$$
\begin{equation*}
a_{i, v}=0 \text { for all v with }|\nu| \leq j \tag{6.16}
\end{equation*}
$$

where $|v|=v_{1}+\ldots+v_{m}$.
If there were an embedding $W_{1} \rightarrow \hat{V}_{m}$ or $W_{1} / \mathbb{R} 1 \rightarrow \hat{V}_{m}$ the Lie algebra $W_{1}$ or $W_{1} / \mathbb{R} 1$ would interit a similar filtration satisfying (6.12) - (6.15). One can now show, essentially by brute force calculations that $W_{1}$ and $W_{1} / \mathbb{R} 1$ do not admit such filtrations. Or much better one observes that (6.12) and (6.14) say that $L_{i} i=0,1,2, \ldots$ is a subalgebra of finite codimension and applies Toby Stafford's result, loc. cit. that $W_{1}$ has no such sub-Lie-algebras.

## 7. PUTTING IT ALL TOGETHER AND CONCLUDING REMARKS

To condude let us spell out the main steps of the argument leading to theorem 2.10 and finish the proof together with some comments as to the generalizability of the various steps.

We start with a stochastic system, in particular the cubic sensor

$$
\begin{equation*}
d x=d w, \quad x(0)=x^{i n}, \quad d y=x^{3} d t+d v \tag{7.1}
\end{equation*}
$$

described more precisely in 2.1 and with a reasonable function $\phi$ of the state of which we want to compute the conditional expectation $\hat{\phi}_{t}$.

The first step now is to show that there exists a pathwise and robust version of $\Phi_{t}$. More precisely it was shown in section 4 that there exist a functional

$$
\begin{equation*}
\delta_{t}^{\phi}(y)=\frac{\Delta_{t}^{\phi}(y)}{\Delta_{t}^{1}(y)}, \quad \Delta_{t}^{\phi}(y)=\left\langle N_{t}(y), \phi\right\rangle \tag{7.2}
\end{equation*}
$$

such that the measures $N_{t}(y)$ depends continuously on the path $y:[0, t] \rightarrow \mathbb{R}$, such that $\Delta_{t}^{1}(y)>0$ all $t>0$, such that the density $n_{t}(y)$ is smooth and such that for $y(t)=y_{t}(\omega)=: y^{\omega}(t)$ a sample path of (7.1) then
(7.3) $\quad \phi_{t}(\omega)=\delta_{t}^{\phi}\left(y^{\omega}\right)$.

From this we also obtained in the case of the cubic sensor that $n_{t}(y)(x)$ as a function of ( $t, x$ ) satisfies the (control version) of the DMZ equation

$$
\begin{equation*}
\frac{\partial}{\partial t} n_{t}(y)(x)=\left(\frac{1}{2} \frac{\partial^{2}}{\partial x^{2}}-\frac{1}{2} x^{6}\right) n_{t}(y)(x)+n_{t}(y)(x) \dot{y}(t) x^{3} \tag{7.4}
\end{equation*}
$$

for piecewise differentiable functions $y:[0, t] \rightarrow \mathbb{R}$. And we showed that the family of densities $n_{t}(y)$, as a function of $t$, is smooth in the sense described in 4.25. Actually a more precise statement is needed, we need smoothness as a function of $t_{1}, \ldots, t_{m}$ if $\dot{y}=u$ with $u$ a bang-bang control of the type $u(t)=\bar{u}_{i} \in \mathbb{R}$ for $t_{1}+\ldots+t_{i-1} \leq t<t_{1}+\ldots+t_{i},\left|u_{i}\right|=1$.

This whole bit is the part of the proof that seems most resistant to
generalization. At present at least this requires reasonable growth bounds on the exponentials occurring in the Kallianpur-Stiebel formula (that is the explicit pathwise expressions for $\left.\Delta_{t}^{\phi}(y)\right)$. In particular let us call a family $\Phi_{t}$ of continuous maps $C_{t} \rightarrow \mathbb{R}$ a path-wise version of $\phi_{t}$, if $\omega \rightarrow \Phi_{t}\left(y^{\omega, t}\right)$, $y^{\omega, t}(s)=y_{s}(\omega), 0 \leq s \leq t$ is a version of $\Phi_{t}$. Then it is not at all clear that path-wise versions exist for arbitrary nonlinear filtering problems.

Now suppose that there exists a smooth finite dimensional filter for $\phi_{t}$. That is a smooth dynamical system

$$
\begin{equation*}
d z=\alpha(z)+\beta(z) d y, \quad \gamma: M \rightarrow \mathbb{R}, \quad z(0)=z^{\text {in }} \tag{7.5}
\end{equation*}
$$

such that if $z_{y}(t)$ denotes the solution of (7.5) then

$$
\begin{equation*}
r\left(z_{y}(t)\right)=\phi_{t}=\delta_{t}^{\phi}(y) \tag{7.6}
\end{equation*}
$$

almost surely. As described in 4.19 above up to a stopping time there also exists a robust pathwise version of the solutions of (7.5) so that $z_{y}(t)$ exists for all continuous $y$ and so that (7.6) holds always. Now let $L_{0}=\frac{1}{2} \frac{d^{2}}{d x^{2}}-\frac{1}{2} x^{6}, L_{3}=x^{3}, L(\bar{u})=L_{0}+u L_{1}$. The next step is to show smoothness of
(7.7) $\quad e^{t_{1} L\left(\bar{u}_{s}\right)} \ldots e^{t_{m} L\left(\bar{u}_{m}\right)} \psi$
for smooth $\psi$ as a function of $t_{1}, \ldots, t_{m}$, and to calculate $\partial^{m} / \partial t_{1} \ldots \partial t_{m}$ of (7.7). The result being formula (4.34).

The next thing is to reinterpret a differential operator on $F$ as a .inear vectorfield $\tilde{L}$ on $F$ by giving meaning to $\tilde{L} \Phi$ for $\Phi$ a functional $F \rightarrow \mathbb{R}$ :or instance a functional of the form $\delta_{t}^{\phi}(y)$.

This permits us to give meaning to expressions like

$$
\begin{equation*}
\left.\frac{\partial^{m}}{\partial t_{1} \ldots \partial t_{m}} \delta_{t}^{\phi}(y)\right|_{t_{1}=\ldots=t_{m}=0} \quad t=t_{1}+\ldots+t_{m} \tag{7.8}
\end{equation*}
$$

for $y \cdot \epsilon C_{t}^{l}$ with $\dot{y}=u$ a bang-bang function. The same operator can be applied to the left handside of $(7.6)$ and as both sides depend smoothly on $t_{1}, \ldots, t_{m}$ there results from (7.6) an equality of the type

$$
\begin{equation*}
\left(A\left(\bar{u}_{m}\right) \ldots A\left(\bar{u}_{1}\right) \gamma\right)(z)=\left(\tilde{L}\left(\bar{u}_{m}\right) \ldots \tilde{L}\left(\bar{u}_{1}\right) \Phi\right)\left(\psi_{z}\right) \tag{7.9}
\end{equation*}
$$

where $z \in M$ and $\psi_{z} \in F$ are corresponding quantities in that they result from feeding in the same control function $\dot{y}(t)$ to the evolution equations for $z$ and $\psi$ respectively.

This relation in turn using some techniques familiar from nonlinear realization theory (essentially restriction to the completely reachable and observable subquotient of $M$ ) then implies that there is a homomorphism of Lie algebras from the Lie algebra $L(\Sigma)$ generated by $L_{0}$ and $L_{1}$ to a Lie algebra of smooth vectorfields. Moreover under the rather inelegant extra assumption that $L(\Sigma)$ contains the operator $\frac{d}{d x} x^{k}$ we showed that $\phi$ must have been constant if this homomorphism of Lie algebras is zero. (Proposition 5.19).

The final part is algebra and shows (i) that $L(\Sigma)=W_{1}$ so that in particular $\frac{d}{d x} x^{k} \in L(\Sigma)$ for all $k=0,1, \ldots$ and (ii) there are no nonzero homomorphisms of Lie algebras $W_{1} \rightarrow V\left(M_{1}\right)$ for $M_{1}$ a smooth finite dimensional manifold. Thus both hypothesis of proposition 5.19 are fulfilled and $\phi$ is a constant. This proves the main theorem 2.10 .

It seems by now clear [Hazewinkel-Marcus, 1981b] that the statement $L(\Sigma)=W_{k}$, $k=\operatorname{dim}$ (state space) will turn out to hold for a great many systems (though anything like a general proof for certain classes of system: is lacking). The system theoretic part of the argument is also quite general The main difficulty of obtaining similar more general results lies thus in generalizing the analytic part or finding suitable subsitutes for establishing the homomorphism principle, perhaps as in [Hijab, 1982].

It should also be stressed that the main theorem 2.10 of this paper only says things about exact filters; it says nothing about approximate filters. On the other hand it seems clear that the Kalman-Bucy filter for $\hat{x}_{t}$ for

$$
\begin{equation*}
d x=d w, \quad d y=x d t+d v \tag{7.10}
\end{equation*}
$$

should for small $\varepsilon$ give reasonable approximate results for

$$
\text { (7.11) } \quad d x=d w, \quad d y=\left(x+\varepsilon x^{3}\right) d t+d v
$$

Yet the estimation Lie algebra of (7.11) is for $\varepsilon \neq 0$ also equal to $W_{1}$ (a somewhat more tedious calculation cf. [Hazewinkel, 1981]) and the arguments of this paper can be repeated word for word (practically) to show that (7.11) does not admit smooth finite dimensional filters (for nonconstant statistics). Positive results that the Kalman-Bucy filter of (7.10) does give an approximation to $\hat{\mathrm{x}}_{\mathrm{t}}$ for (7.11) are contained in loc. cit. [Sussmann, 1982], and [Blankenship - Liu - Marcus, 1983].

It is possible that results on approximate filters can be obtained by considering $L(\Sigma)$ not as a bare Lie algebra but as a Lie algebra with two distinguished generations $L_{0}, L_{1}$ which permits us to consider also the Lie algebra $L_{s}(\Sigma)$ generated by $s L_{0}, s L_{1}$ (where $s$ is an extra variable) and to consider statements like $L_{s}(\Sigma)$ is close to $L_{s}\left(\Sigma^{\prime}\right)$ module $s^{t}$.

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