# A General Notion of Shears, and Applications 

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## 1. Introduction

In this paper we introduce a generalization of the notions of shears and overshears to arbitrary complex manifolds. The concept is very simple, but it is useful in the study of complex manifolds having very large automorphism groups. We shall explore some of the consequences of this concept in connection with the density property, which we now recall.

In [V1] we introduced the notion of complex manifolds with the density property. Recall that a complex manifold $M$ has the density property if the Lie subalgebra of $\mathcal{X}_{\mathcal{O}}(M)$ generated by the complete vector fields on $M$ is a dense subalgebra. More generally, a Lie subalgebra $\mathfrak{g} \subset \mathcal{X}_{\mathcal{O}}(M)$ is said to have the density property if the complete vector fields in $\mathfrak{g}$ generate a dense subalgebra of $\mathfrak{g}$. (So $M$ has the density property if and only if $\mathcal{X}_{\mathcal{O}}(M)$ has the density property.) Another important case occurs when $M$ has a nonvanishing holomorphic $n$-form $\left(n=\operatorname{dim}_{\mathbb{C}} M\right)$, that is, a holomorphic volume element $\omega$. We say that $(M, \omega)$ has the volume density property if the Lie algebra $\mathcal{X}_{\mathcal{O}}(M, \omega):=\left\{X \in \mathcal{X}_{\mathcal{O}}(M) \mid L_{X} \omega=0\right\}$ has the density property. Andersén $[\mathrm{A}]$ proved that $\left(\mathbb{C}^{n}, d z_{1} \wedge \cdots \wedge d z_{n}\right)$ has the volume density property, and then Andersén and Lempert [AL] proved that $\mathbb{C}^{n}$ has the density property. The author showed that for every complex Lie group $G,(G \times \mathbb{C}, \omega)$ has the volume density property, where $\omega$ is the unique (up to constant multiple) left (or right) invariant holomorphic volume element on $G \times \mathbb{C}$, and that if $G$ is a Stein Lie group, then $G \times \mathbb{C}$ has the density property. The author also produced several examples of Lie algebras of vector fields with the density property.

In [V2] we used jets to explore the complex structure of (mostly Stein) complex manifolds with the density property. It was shown, among other things, that Stein manifolds with the density property admit open subsets biholomorphic to $\mathbb{C}^{n}$ and have interesting properties with respect to their embedded submanifolds. Some of the results were known for $\mathbb{C}^{n}$ through works of Buzzard, Fornæss, Forstnerič, Globevnik, Rosay, Stensønes, and others.

With the usefulness of the density property already established in the literature, some sort of classification or fine structure theorem is very desirable. Such a result seems at the moment very far off, owing in part to the lack of examples. The main
theorems of this paper, which we now state, give many new examples of the density property; more importantly, the proofs establish techniques that can be used to construct other examples. We shall pursue this in future work.

Theorem 1. Let $M^{2}:=\mathbb{C}^{2} \backslash\{x y=1\}$ and $\omega:=(x y-1)^{-1} d x \wedge d y$. Then $\left(M^{2}, \omega\right)$ has the volume density property.

The study of the space $M^{2}$ was inspired by discussions with Rosay several years ago. This space is important because it is another instance of the mysterious prephenomenon (we say "pre" because there are no proofs that it exists) of a holomorphic volume element that is preserved by every holomorphic automorphism.

In the next result, we study a complex Lie group that is not of the form $G \times \mathbb{C}$. There is, as of yet, no general theory here, so we focus on one example.

Theorem 2. The complex Lie group $\mathrm{Sl}(2, \mathbb{C}):=\left\{(a, b, c, d) \in \mathbb{C}^{4} \mid a d-b c=\right.$ 1\} has the density and volume density property.

Next we introduce a new class of complex manifolds with holomorphic volume element called EMV manifolds. These spaces are generalizations of complex Lie groups, but also of certain complex homogeneous spaces. Roughly speaking, they have the property that all holomorphic vector fields on them can be approximately written as finite sums of the form $\sum f_{j} X_{j}$, where $f_{j}$ are any holomorphic functions, and $X_{j}$ are divergence zero completely generated holomorphic vector fields (see Section 2).

Theorem 3. Let $(M, \omega)$ be an EMV manifold. Then $(M \times \mathbb{C}, \omega \wedge d z)$ has the volume density property. If $M$ is moreover an open subset of a Stein manifold, then $M \times \mathbb{C}$ has the density property.

As already suggested, the key tool used in the proofs of these theorems is a generalization to arbitrary complex manifolds of the notion of shears and overshears. This tool may have some independent interest as well. The idea is quite simple: Given a $\mathbb{C}$-complete holomorphic vector field $X$ in a complex manifold $M$, one tries to produce new complete vector fields of the form $f \cdot X$, with $f \in \mathcal{O}(M)$. We establish necessary and sufficient conditions on such $f$, and these conditions define in a natural way function spaces associated to $X$. We then prove theorems to the effect that the structure of these function spaces depends on the intrinsic and extrinsic geometry of the orbits of $X$.

The organization of the paper is as follows. In Section 2 we briefly recall some basic definitions in the theory of ordinary differential equations and volume geometry, taking the opportunity to establish notation. In Section 3 we introduce and develop general shears and overshears. In part, our results here explain why it was easiest to prove the density property for spaces of the form $G \times \mathbb{C}$. In Section 4 we prove Theorem 1, and in Section 5 we prove Theorem 2; the proofs are rather combinatorial in nature. In Section 6 we introduce EM and EMV spaces, and we prove Theorem 3 as well as some related results. Finally, in Section 7 we state a question which naturally arises in the course of the paper, giving an example of a
complex manifold which may or may not have the (volume) density property but for which the combinatorial methods of Sections 4 and 5 become too cumbersome to carry out.

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## 2. Some Preliminaries

In this section we recall a few basic concepts and establish the notation used below.
A holomorphic vector field $X$ is a holomorphic section of $T^{1,0} M$, the holomorphic part of the complexified tangent bundle. Since there is a natural identification of $T^{1,0} M$ with the real tangent bundle $T M$, we can identify $X$ with a real vector field, which we still denote by $X$. This vector field has a flow $\varphi_{X}$, which is a map defined on an open subset of $M \times \mathbb{R}$ containing $M \times\{0\}$ as follows: for $(x, t) \in$ $M \times \mathbb{R}, \varphi_{X}^{t}(x)$ is the point $c(t) \in M$, where $c: I \subset \mathbb{R} \rightarrow M$ is the maximal solution of the initial value problem

$$
\frac{d c}{d t}=X(c), \quad c(0)=x
$$

Moreover, $\varphi_{X}^{t}$ is holomorphic for each $t$. We denote the set of holomorphic vector fields on $M$ by $\mathcal{X}_{\mathcal{O}}(M)$.

A holomorphic vector field $X$ is called complete if $\varphi_{X}$ is defined on all of $M \times \mathbb{R}$. In this case $\left\{\varphi_{X}^{t} \mid t \in \mathbb{R}\right\}$ is a 1-parameter group of automorphisms of $M$.

A vector field $X$ is called $\mathbb{C}$-complete if both $X$ and $i X$ are complete. Let $\psi^{s+i t}(x):=\varphi_{X}^{s} \circ \varphi_{i X}^{t}(x)$. One checks that, since $[X, i X]=0$ for all holomorphic vector fields, $\left\{\psi^{\zeta} \mid \zeta \in \mathbb{C}\right\}$ defines a complex 1-parameter group of automorphisms which is holomorphic in $\zeta$, that is, a holomorphic $\mathbb{C}$-action. In this paper we shall use the term complete to mean $\mathbb{C}$-complete.

The set $\mathcal{X}_{\mathcal{O}}(M)$ of all holomorphic vector fields on $M$ is equipped with a bracket (or commutator) operation, $[X, Y]=X Y-Y X$, which makes it into a Lie algebra. Given any Lie algebra $\mathfrak{g}$ of holomorphic vector fields, we can consider the Lie subalgebra $\mathfrak{g}^{\prime}$ of $\mathfrak{g}$ generated by the complete vector fields in $\mathfrak{g}$. Any $X \in \overline{\mathfrak{g}^{\prime}}$ is said to be $\mathfrak{g}$-completely generated. If $\mathfrak{g}=\mathcal{X}_{\mathcal{O}}(M)$, we omit reference to the Lie algebra. If $\mathfrak{g}=\mathcal{X}_{\mathcal{O}}(M, \omega)$ (see below), we say that $X \in \overline{\mathfrak{g}^{\prime}}$ is divergence zero completely generated.

Let us now suppose that $M$ admits a nowhere vanishing holomorphic $n$-form $\omega$, where $n=\operatorname{dim}_{\mathbb{C}} M$. We call such a form a holomorphic volume element. Given a holomorphic volume element $\omega$, we can define a map $\operatorname{div}_{\omega}: \mathcal{X}_{\mathcal{O}}(M) \rightarrow \mathcal{O}(M)$ by

$$
\operatorname{div}_{\omega}(X)=\frac{L_{X} \omega}{\omega}
$$

where $L_{X}$ is the Lie derivative of $X$ :

$$
L_{X} \alpha=\left.\frac{d}{d t}\right|_{t=0}\left(\varphi_{X}^{t}\right)^{*} \alpha
$$

Since $L_{[X, Y]}=L_{X} L_{Y}-L_{Y} L_{X}$, one easily shows that

$$
\operatorname{div}_{\omega}([X, Y])=X \operatorname{div}_{\omega} Y-Y \operatorname{div}_{\omega} X
$$

Another useful formula, due to H. Cartan, is

$$
\operatorname{div}_{\omega}(X)=\frac{d\left(i_{X} \omega\right)}{\omega}
$$

where $i_{X}$ is contraction with respect to $X$.
Finally, we denote the kernel of $\operatorname{div}_{\omega}$ by $\mathcal{X}_{\mathcal{O}}(M, \omega)$ and call $X \in \mathcal{X}_{\mathcal{O}}(M, \omega)$ a divergence zero vector field.

## 3. General Shears and Overshears

## Basic Propositions and the Definition

Let $X \in \mathcal{X}_{\mathcal{O}}(M)$. We define

$$
I^{j}(X)=I_{\mathcal{O}}^{j}(X):=\left\{f \in \mathcal{O}(M) \mid X^{j} f=0\right\}
$$

If $f \in I^{1}(X)$ (resp. $I^{2}(X)$ ), we say $f$ is a first (resp. second) integral of $X$. The following proposition is immediate.

Proposition 3.1. Let $X$ be a holomorphic vector field with (local) flow $g_{X}^{t}$. Then $f \in I^{1}(X)$ (resp. $\left.I^{2}(X)\right)$ if and only if (where defined)

$$
f \circ g_{X}^{t}=f \quad\left(\text { resp. } f \circ g_{X}^{t}=f+t X f\right)
$$

Although first integrals have been studied extensively in the past, second integrals seem not to have been looked at. However, in the holomorphic category it is natural to study first and second integrals because of the following fundamental proposition.

Proposition 3.2. If $X \in \mathcal{X}_{\mathcal{O}}(M)$ is $\mathbb{C}$-complete and $f \in \mathcal{O}(M)$, then $f X$ is $\mathbb{C}$ complete if and only if $f \in I^{2}(X)$.

Proof. If $X$ vanishes at some $p \in M$, then so does $f X$, so the integral curve of $f X$ through $p$ is defined (and constant) for all $t \in \mathbb{C}$. Suppose now that $X(p) \neq$ 0 . Let $h_{p}: \mathbb{C} \rightarrow R_{p}(X)$ be the integral curve of $X$ through $p$. Here, $R_{p}(X)$ is the orbit of $X$ through $p$. Then

$$
h_{p}^{*}(X)(t)=\partial_{t},
$$

and $h_{p}$ is a covering map. Since $f X$ is tangent to the orbits of $X, h_{p}^{*}(f X)$ is a well-defined vector field on $\mathbb{C}$. Precisely,

$$
h_{p}^{*}(f X)(t)=f \circ h_{p}(t) \partial_{t}
$$

It follows that the integral curve of $f X$ through $p$ is defined for all time if and only if $f \circ h_{p}(t)$ is an affine linear function of $t$. This holds for all $p$ in $M \backslash$ $\{X=0\}$ if and only if $f \in I^{2}(X)$.

Proposition 3.2 is a purely holomorphic result. Note that, in general, multiplying a (real) vector field by any bounded function preserves completeness.

Example 3.3. Let $\mathbb{C}^{n}=\mathbb{C} \times \mathbb{C}^{n-1}(n \geq 2)$ have coordinates $z=\left(z_{1}, z^{\prime}\right)$. Consider the vector field $\partial_{z_{1}}$ on $\mathbb{C}^{n}$. Then $f(z) \partial_{z_{1}}$ is complete if and only if $f(z)=g\left(z^{\prime}\right)+h\left(z^{\prime}\right) z_{1}$. Vector fields of the form $f\left(z^{\prime}\right) \partial_{z_{1}}$ are called shear fields; those of the form $f\left(z^{\prime}\right) z_{1} \partial_{z_{1}}$ are called overshear fields. These vector fields have played a fundamental role in the study of automorphisms of $\mathbb{C}^{n}$, as the set of all time-1 maps of these vector fields generates a dense subgroup of Aut $\left(\mathbb{C}^{n}\right)$ [A; AL].

Definition 3.4. Let $X$ be a complete holomorphic vector field on a complex manifold $M$. An $X$-shear (resp. $X$-overshear) field on $M$ is a vector field of the form $f \cdot X$, with $f \in I^{1}(X)\left(\operatorname{resp} . I^{2}(X)\right)$.

## Second Integrals

To find first integrals of a complete vector field $X$, it is well known that the orbits of $X$ must have particularly nice behavior. Since $X$ maps $I^{2}(X)$ to $I^{1}(X)$, we can expect that second integrals are somehow more rare than first integrals. We will show that this is indeed the case.

One can phrase the problem of finding second integrals (i.e., solving the secondorder PDE $X^{2} f=0$ ) as an inhomogeneous first-order PDE with conditions on the forcing term:

$$
X v=\varphi \quad \text { with } \varphi \in I^{1}(X)
$$

The most optimistic situation occurs when we can solve the equation $X u=1$. In this case, we can write $f \in I^{2}(X)$ as

$$
f=u X f+(f-u X f)
$$

which shows that $I^{2}(X)=I^{1}(X)+u I^{1}(X)$. We shall see, however, that $X u=1$ does not always have a solution.

To get a good idea of when $I^{2}(X)$ is "large", it is convenient to use the language of ideals. Let $J_{X}:=X\left(I^{2}(X)\right) \subset I^{1}(X)$. $J_{X}$ is an ideal in $I^{1}(X)$, since $\varphi X f=X(\varphi f) \in J_{X}$ for $\varphi \in I^{1}(X)$ and $f \in I^{2}(X)$. Being able to solve $X u=1$ is equivalent to saying that $J_{X}=I^{1}(X)$. Hence $I^{2}(X)$ is "large" when $I^{1}(X)$ is large and the quotient ring $I^{1}(X) / J_{X}$ is "small"-for example, finitely generated or trivial.

It is interesting that the size of the quotient $I^{1}(X) / J_{X}$ is intimately tied up with the complex geometry of the orbit space of $X$. Our first result is the following.

Theorem 3.5. Let $X \in \mathcal{X}_{\mathcal{O}}(M)$ be complete, let $f \in I^{2}(X)$, and set $N:=$ $M \backslash\{X f=0\}$. Then
(1) $N / X$ is a complex manifold,
(2) $\pi: N \rightarrow N / X$ is a holomorphic submersion,
(3) $\pi \times f: N \rightarrow(N / X) \times \mathbb{C}$ is a biholomorphic map, and

$$
\begin{equation*}
(\pi \times f)_{*}(f X)\left(R_{p}(X), \lambda\right)=\psi\left(R_{p}(X)\right) \lambda \partial_{\lambda} \tag{4}
\end{equation*}
$$

for some $\psi \in \mathcal{O}(N / X)$.
Proof. Let $u:=(1 / X f) f$. Then $X u=1$, and $u \circ g^{t}(p)=u(p)+t$. Note also that $\left.X\right|_{N}$ is complete, since $\{X f=0\}$ is a union of orbits.
(1) The manifold $N / X$ can be identified with the level set $u^{-1}(0)$ via the map

$$
\xi: N / X \rightarrow u^{-1}(0) ; R_{p}(X) \mapsto R_{p}(X) \cap u^{-1}(0)
$$

First, if $p \in N$ then $g_{X}^{-u(p)}(p) \in u^{-1}(0)$, so that no orbit has empty intersection with $u^{-1}(0)$. Hence $\xi$ is well-defined, at least as a set-valued function. Next, note that $\xi$ is single-valued. Indeed, if $R_{p}(X) \cap u^{-1}(0)$ contains $p_{1}$ and $p_{2}$, then $p_{1}=$ $g_{X}^{t_{1}}(p)$ and $p_{2}=g_{X}^{t_{2}}(p)$. But since $u\left(p_{1}\right)=u\left(p_{2}\right)$ and $u \circ g^{t}(p)=u(p)+t$, we see that $t_{1}=t_{2}$ and hence that $p_{1}=p_{2}$. Next, $\xi$ is $1-1$ because orbits of vector fields never intersect. Finally, $\xi$ is clearly surjective. To finish (1), note that since $d u(X)=1, d u$ never vanishes on $N$. Hence $u^{-1}(0)$ is a complex manifold, which we henceforth identify with $N / X$ via $\xi$.
(2) Observe that the canonical projection $\pi: N \rightarrow u^{-1}(0)$ is given by $\pi(p)=$ $g_{X}^{-u(p)}(p)$. Note also that $\left.\pi\right|_{u^{-1}(t)}: u^{-1}(t) \rightarrow u^{-1}(0)$ is a biholomorphic map; $\left.\pi\right|_{u^{-1}(t)}=g_{X}^{-t}$. Hence $\pi$ is a submersion.
(3) Define $\tau: N \times \mathbb{C} \rightarrow \mathbb{C}$ and $G: N \times \mathbb{C} \rightarrow N$ by

$$
\tau(p, \lambda):=\frac{\lambda-f(p)}{X f(p)} \quad \text { and } \quad G(p, \lambda):=g_{X}^{\tau(p, \lambda)}(p)
$$

Then, since $f \in I^{2}(X)$ (and hence $X f \in I^{1}(X)$ ), Proposition 3.1 gives that $\tau\left(g_{X}^{t}(p), \lambda\right)=\tau(p, \lambda)-t$ and hence that

$$
G\left(g_{X}^{t}(p), \lambda\right)=g_{X}^{\tau(p, \lambda)-t} \circ g_{X}^{t}(p)=G(p, \lambda) .
$$

Thus $G$ defines a holomorphic map $H:(N / X) \times \mathbb{C} \rightarrow N$ by

$$
H\left(R_{p}(X), \lambda\right):=G(p, \lambda)
$$

Now

$$
\begin{aligned}
\pi \times f \circ H\left(R_{p}(X), \lambda\right) & =\pi \times f\left(g_{X}^{\tau(p, \lambda)}(p)\right. \\
& =\left(R_{p}(X), f(p)+\tau(p, \lambda) X f(p)\right) \\
& =\left(R_{p}(X), \lambda\right)
\end{aligned}
$$

and

$$
\begin{aligned}
H \circ \pi \times f(p) & =H\left(\left(R_{p}(X), f(p)\right)\right. \\
& =g_{X}^{\tau(p, f(p))}(p) \\
& =p .
\end{aligned}
$$

Thus $H=(\pi \times f)^{-1}$ and hence $\pi \times f$ is a biholomorphic map.
(4) We have

$$
\begin{aligned}
(\pi \times f)_{*}(X) & =\left.\frac{d}{d t}\right|_{t=0} \pi \times f \circ g_{X}^{t} \circ H\left(R_{p}(X), \lambda\right) \\
& =\left.\frac{d}{d t}\right|_{t=0} \pi \times f \circ g_{X}^{t} \circ g_{X}^{\tau(p, \lambda)}(p) \\
& =\left.\frac{d}{d t}\right|_{t=0}\left(R_{p}(X), f(p)+\tau(p, \lambda) X f(p)+t X f(p)\right) \\
& =X f(p) \partial_{\lambda},
\end{aligned}
$$

so now

$$
\begin{aligned}
\left((\pi \times f)_{*}(f X)\right)\left(R_{p}(X), \lambda\right) & =\left(H^{*} f\right)\left(R_{p}(X), \lambda\right) \cdot(\pi \times f)_{*}(X)\left(R_{p}(X), \lambda\right) \\
& =\lambda X f(x) \partial_{\lambda} .
\end{aligned}
$$

Taking $\psi\left(R_{p}(X)\right)=X f(p)$ finishes the proof.
As a corollary, we obtain the following proposition.
Proposition 3.6. Let $X \in \mathcal{X}_{\mathcal{O}}(M)$ be complete, and define

$$
\Sigma_{X, M}:=\bigcap_{f \in I^{2}(X)}\{X f=0\}, \quad N_{X, M}:=M \backslash \Sigma_{X, M}
$$

(Note that $N_{X, M}$ is an open subset of $M$, which is either empty or dense.) Then for each $p \in N_{X, M}, R_{p}(X)$ is biholomorphic to $\mathbb{C}$. In particular, if $X$ has a nontrivial second integral, then almost every orbit of $X$ is biholomorphic to $\mathbb{C}$.

Suppose we can solve $X v=\varphi \in I^{1}(X)$. Then Theorem 3.5 tells us that $N / X$ is a complex manifold, and $N(=M \backslash\{\varphi=0\})$ is biholomorphic to $N / X \times \mathbb{C}$. It follows that if $M$ is Stein then $N / X$ is itself Stein (since $N$ is Stein). In the case where $\varphi \equiv 1$, the converse is also true.

Theorem 3.7. Let $X \in \mathcal{X}_{\mathcal{O}}(M)$ be a complete vector field whose orbits all are biholomorphic to $\mathbb{C}$. Suppose $M / X$ is a complex manifold and $\pi: M \rightarrow M / X$ is a holomorphic map. If $M / X$ is Stein, then $X u=1$ has a solution.

Proof. If $M / X$ is a (differentiable) manifold and $\pi$ is smooth, then $\pi$ is a submersion and thus the bundle $\pi: M \rightarrow M / X$ is locally trivial. Furthermore, it is possible to select local trivializations $\left\{\varphi_{j}: \pi^{-1}\left(U_{j}\right) \rightarrow U_{j} \times \mathbb{C}\right\}$ such that $\left(\varphi_{j}\right)_{*} X=$ $\partial_{\lambda}$ for all $j$. Indeed, let $\sigma_{j}$ be a local section of $\pi: M \rightarrow M / X$ over $U_{j}$. For each $x \in \pi^{-1}\left(U_{j}\right)$, define $\lambda=\lambda(x)$ to be the unique complex number for which $g_{X}^{\lambda}\left(\sigma_{j} \circ \pi(x)\right)=x$. The dependence of $\lambda$ on $x$ is holomorphic because of the holomorphic dependence of the flow on initial conditions. Set $\varphi_{j}(x):=(\pi(x), \lambda(x))$. Note that $\varphi_{j} \circ g_{X}^{s}(x)=(\pi(x), s+\lambda(x))$ and so

$$
\left(\varphi_{j}\right)_{*} X(x)=\left.\frac{d}{d s}\right|_{s=0}(\pi(x), s+\lambda(x))=\partial_{\lambda} .
$$

Now, since the fibers of our holomorphic bundle are $\mathbb{C}$, the bundle must be an affine bundle; hence the transition functions $\varphi_{j k}(\pi(x)) t:=\operatorname{pr}_{\lambda} \circ \varphi_{j} \circ \varphi_{k}^{-1}(x, t)$ (where $\mathrm{pr}_{\lambda}$ is the projection to the second factor) satisfy

$$
\varphi_{j k}(\pi(x)) t=f_{j k}(\pi(x)) t+g_{j k}(\pi(x))
$$

Moreover, because of the way the $\varphi_{j}$ were chosen, $f_{j k}(\pi(x)) \equiv 1$ for all $j, k$. Indeed,

$$
f_{j k}(\pi(x))=\frac{\partial}{\partial t} \varphi_{j k}(\pi(x)) t=\operatorname{pr}_{\lambda *}\left(\varphi_{j}\right)_{*}\left(\varphi_{k}^{-1}\right)_{*} \partial_{\lambda} \equiv 1
$$

Next, writing out the identity

$$
\varphi_{j k} \circ \varphi_{k l} \circ \varphi_{l j}=\mathrm{id}
$$

shows that $\left\{g_{j k}\right\}$ is a 1-cocycle on $M / X$ (i.e., Cousin-1 data). Since $M / X$ is Stein, $g_{j k}=g_{j}-g_{k}$. One checks easily that $\left\{g_{k}\right\}$ is a section of $\pi: M \rightarrow M / X$. It follows that $\pi: M \rightarrow M / X$ is actually a line bundle, since we can use the section $\left\{g_{k}\right\}$ as an origin for each fiber. Precisely, we can define the transition functions

$$
P_{j k}(\pi(x)) v:=\varphi_{j k}(\pi(x))\left(v+g_{k}(\pi(x))\right)-g_{j}(\pi(x)) .
$$

Then
$P_{j k}(\pi(x)) v=f_{j k}(\pi(x))(v)+\varphi_{j k}(\pi(x))\left(g_{k}(\pi(x))\right)-g_{j}(\pi(x))=f_{j k}(\pi(x))(v)$
so that, since $f_{j k} \equiv 1, \pi: M \rightarrow M / X$ is trivial. We now define (in the usual way) the global trivialization $F: M \rightarrow(M / X) \times \mathbb{C}$ by $F:=\pi \times \psi$, where

$$
\psi(x)=\operatorname{pr}_{\lambda} \circ \varphi_{j}(x)-g_{j}(\pi(x)) \quad \text { for } x \in \pi^{-1}\left(U_{j}\right)
$$

The function $\psi$ is well-defined, since for $x \in U_{j} \cap U_{k}$ we have

$$
\begin{aligned}
\operatorname{pr}_{\lambda} \circ \varphi_{j}(x)-g_{j}(\pi(x)) & =\operatorname{pr}_{\lambda} \circ \varphi_{j} \circ \varphi_{k}^{-1}\left(\varphi_{k}(x)\right)-g_{j}(\pi(x)) \\
& =\varphi_{j k}\left(\pi\left(\varphi_{k}(x)\right) t-g_{j}(\pi(x)) \quad \text { where } t=\operatorname{pr}_{\lambda}\left(\varphi_{k}(x)\right)\right. \\
& =t+g_{j k}(\pi(x))-g_{j}(\pi(x)) \\
& =t-g_{k}(\pi(x)) \\
& =\operatorname{pr}_{\lambda} \circ \varphi_{k}(x)-g_{k}(\pi(x))
\end{aligned}
$$

It follows that

$$
F_{*} X=\partial_{\lambda}
$$

Setting $u\left(F^{-1}(\pi(x), \lambda)\right)=\lambda$, we see that $X u=X\left(F^{*}\left(F_{*} u\right)\right)=\left(F_{*} X\right)\left(F_{*} u\right)=1$, as required.

Remarks. (1) A more careful look at the proof shows that one does not need $M / X$ to be Stein, but only that $H^{1}(M / X, \mathcal{O})=0$.
(2) Theorems 3.5 and 3.7 explain in part why it was so much easier to prove density theorems for spaces of the form $M \times \mathbb{C}$.

Example 3.8. Let

$$
X(x)=a \partial_{b}+c \partial_{d} \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C})), \quad x=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in \operatorname{Sl}(2, \mathbb{C}) .
$$

Note that $X$ is a left invariant vector field on $\operatorname{Sl}(2, \mathbb{C})$ whose orbits are closed and biholomorphic to $\mathbb{C}$. Hence $\mathrm{Sl}(2, \mathbb{C}) / X$ is a complex manifold. Nevertheless, the equation $X u=1$ has no global holomorphic solutions. Indeed, $\mathrm{Sl}(2, \mathbb{C})$ is homotopy equivalent (by the Gram-Schmidt algorithm) to $S U(2) \cong S^{3}$, which is a cell complex of dimension 3. It follows that $\operatorname{Sl}(2, \mathbb{C})$ is not biholomorphic to $B \times \mathbb{C}$, for then $B$ would be a Stein 2 -fold with 3 -dimensional cells, a contradiction. (It is interesting to note, however, that $u=(\bar{a} b+\bar{c} d) /\left(|a|^{2}+|c|^{2}\right)$ is a real analytic solution and that the $\mathbb{C}$-fibration $\mathrm{Sl}(2, \mathbb{C}) \rightarrow \mathrm{Sl}(2, \mathbb{C}) / X$ is real-analytically trivial.)

## 4. $\left(M^{2}, \omega\right)$

Recall that we define

$$
M^{2}=\mathbb{C}^{2} \backslash\{x y=1\} \quad \text { and } \quad \omega=\frac{1}{x y-1} d x \wedge d y
$$

In this section we prove Theorem 1.

## Notation and Facts

It will be convenient to write $z=x y-1$. As we mentioned before, $M^{2}$ admits two everywhere independent complete vector fields,

$$
X(x, y):=z \partial_{y} \quad \text { and } \quad Y(x, y):=z \partial_{x} .
$$

Since $z$ does not vanish on $M^{2}$, it is clear that every holomorphic vector field on $M^{2}$ is of the form $f X+g Y$ for some $f, g \in \mathcal{O}\left(M^{2}\right)$. We note also that

$$
H(x, y)=x \partial_{x}-y \partial_{y}
$$

is a complete holomorphic vector field with zero divergence. One can integrate $X, Y$, and $H$ to see that every orbit of $X$ and $Y$ is biholomorphic to $\mathbb{C}^{*}$ and that this is also the case for every orbit of $H$, except for its single fixed point at the origin of $\mathbb{C}^{2}$. Hence $X, Y$, and $H$ have no nontrivial second integrals. The following facts are easily computed:

$$
\begin{gathered}
{[H, X]=X, \quad[H, Y]=-Y, \quad[X, Y]=z H,} \\
x Y-y X=z H, \\
X x=0, \quad X y=z, \quad X z=x z \\
Y x=z, \quad Y y=0, \quad Y z=y z \\
H x=x, \quad H y=-y, \quad H z=0
\end{gathered}
$$

Lemma 4.1. Every $\varphi \in \mathcal{O}\left(M^{2}\right)$ is of the form

$$
\varphi(x, y)=f(x, y, z)
$$

for some $f \in \mathcal{O}\left(\mathbb{C}^{2} \times \mathbb{C}^{*}\right)$.
Proof. The mapping $j:(x, y) \mapsto(x, y, z)$ gives a proper holomorphic embedding of $M^{2}$ into $\mathbb{C}^{2} \times \mathbb{C}^{*}$. It is thus a standard fact (Theorem A) that $\mathcal{O}\left(M^{2}\right)=$ $\left.\mathcal{O}\left(\mathbb{C}^{2} \times \mathbb{C}^{*}\right)\right|_{M^{2}}$.

Thus the Laurent polynomials

$$
\sum_{j, l \geq 0} c_{j l} x^{j} z^{-l}+\sum_{k, l \geq 0} d_{k l} y^{k} z^{-l}+\sum_{l \geq 0} e_{l} z^{-l}+\sum_{k, l \geq 0} f_{j k} x^{j} y^{k}
$$

are dense in $\mathcal{O}\left(M^{2}\right)$. We shall call such Laurent polynomials reduced.

## The Key Lemmas

Lemma 4.2. Let $j$ and $k$ be nonnegative integers. Then, for some polynomial $p(x, y)$, there is a divergence zero completely generated vector field of the form

$$
x^{j} y^{k} X+p(x, y) Y
$$

Proof. Since $X x=0, x^{j} X$ is complete, which proves the claim for $k=0$. Note next that, since $Y y=0, y^{l} Y$ is complete, and hence (as a computation shows)

$$
\left[y^{k} Y, x^{j} X\right]=(j+1) x^{j} y^{k+1} X-j x^{j-1} y^{k} X+p_{1}(x, y) Y
$$

is divergence zero completely generated. The result follows by induction on $k$.
This lemma has a corollary which is of independent interest. Let $\mathfrak{g}$ denote the Lie algebra of all holomorphic vector fields of $\mathbb{C}^{2}$ that vanish on $\{x y=1\}$ and have $\omega$-divergence zero.

Corollary 4.3. The Lie algebra $\mathfrak{g}$ has the density property.
Proof. Note first that the set of divergence zero vector fields of the form $p(x, y) X+q(x, y) Y$ is dense in $\mathfrak{g}$ for polynomials $p$ and $q$. Let $V$ be one such vector field. By Lemma 4.2 there exists another such vector field $W$, which is completely generated, such that $V-W=p_{1}(x, y) Y$. But since $0=\operatorname{div}(V-W)=$ $Y\left(p_{1}\right), V-W$ is complete. Thus $V=W+(V-W)$ is completely generated, as desired.

The following identities are obtained by simple computations, the last two most easily proved using the commutation relations given above. We omit the details.

$$
\begin{aligned}
-z^{-l} H & =y z^{-l} X+(*) Y, \\
{\left[z^{-l} H, y^{k} Y\right] } & =l y^{k+2} z^{-(l+1)} X+(*) Y, \\
{\left[x^{j} X, z^{-l} H\right] } & =l x^{j} z^{-(l+1)} X+(l-j-1) x^{j} z^{-l} X+(*) Y .
\end{aligned}
$$

Here and in what follows, the symbol $(*)$ means a polynomial in $x, y$, and $1 / z$. Using the first identity, we have the following lemma.

Lemma 4.4. For each $l \geq 1$, there exists a complete divergence zero vector field of the form

$$
y z^{-l} X+(*) Y
$$

Using the second identity, we have the following.
Lemma 4.5. For each $l \geq 2$ and $k \geq 0$, there exists a divergence zero completely generated vector field of the form

$$
y^{k+2} z^{-l} X+(*) Y .
$$

Using the third identity, by induction we have the following lemma.
Lemma 4.6. For each $l \geq 2$ and $j \geq 0$, there exists a polynomial $p(x)$ and $a$ divergence zero completely generated vector field of the form

$$
y^{j} z^{-l} X+p(x) z^{-1} X+(*) Y .
$$

Lemma 4.7. Suppose that $p$ and $q$ are polynomials in one variable, that $g \in$ $\mathcal{O}\left(M^{2}\right)$, and that

$$
V(x, y)=\frac{p(x)+y q(y)}{z} X+g(x, y) Y
$$

is a divergence zero vector field. Then $p=0$ and $q$ is constant.
Proof. The vanishing divergence of $V$ is equivalent to the closedness of the holomorphic 1 form $\theta=i_{V} \omega$. An easy computation shows that

$$
\theta=-\frac{p(x)+y q(y)}{z} d x+g(x, y) d y
$$

It follows from Stokes's theorem that, if $\Omega$ is a smooth 2-manifold with boundary, then

$$
\int_{\partial \Omega} \theta=0 .
$$

For $y \in \mathbb{C}^{*}$, let $\gamma_{y}:[0,2 \pi] \rightarrow M^{2}$ be defined by

$$
\gamma_{y}(t)=\left(\left(e^{i t}+1\right) y, y^{-1}\right)
$$

Note that

$$
\begin{aligned}
\int_{\gamma_{y}} \theta & =\int_{0}^{2 \pi} \frac{p\left(\left(1+e^{i t}\right) y\right)+(1 / y) q(1 / y)}{e^{i t}} i y e^{i t} d t \\
& =2 \pi i(y p(y)+q(1 / y))
\end{aligned}
$$

Fix $y_{0}$ and $y_{1}$ in $\mathbb{C}^{*}$, and let $\beta:[0,1] \rightarrow \mathbb{C}^{*}$ be any smooth curve with $\beta(0)=y_{0}$ and $\beta(1)=y_{1}$. Then

$$
\Omega_{y_{0}, y_{1}}:=\left\{\gamma_{\beta(s)}(t) \mid(t, s) \in[0,1] \times[0,2 \pi]\right\}
$$

is a smooth cylinder in $M^{2}$, and

$$
\partial \Omega=\gamma_{y_{0}} \cup \gamma_{y_{1}} .
$$

Since $y_{0}, y_{1}$ were arbitrary, it follows that the Laurent polynomial $y p(y)+q(1 / y)$ is constant and hence that $p=0$ and $q$ is constant. This completes the proof.

Proof of Theorem 2. Let $V=f X+g Y$ be a holomorphic vector field with $f$ and $g$ reduced (see the remark following Lemma 4.1) Laurent polynomials. By Lemmas 4.2, 4.5, and 4.6, there exists a divergence zero completely generated vector field $W_{1}$ such that $V-W_{1}=((p(x)+y q(y)) / z) X+(*) Y$. According to Lemma 4.7, $p=0$ and $q$ is constant. Hence, by Lemma 4.4, there is a complete vector field $W_{2}$ such that $V-W_{1}-W_{2}=h(x, y) Y$ for some $h \in \mathcal{O}\left(M^{2}\right)$. But since $Y h=$ $\operatorname{div}(h Y)=0, V-W_{1}-W_{2}$ is complete. Hence

$$
V=W_{1}+W_{2}+\left(V-W_{1}-W_{2}\right)
$$

is divergence zero completely generated, as desired.
As mentioned in Section 1, it is not known whether there exists a single automorphism $f$ of $M^{2}$ such that $f^{*} \omega \neq \pm \omega$. However, this difficulty is immediately lifted by "stabilizing" $M^{2}$. Theorem 2 and the main result I. 3 in [V1] imply the following.

Corollary 4.8. $\quad M^{2} \times \mathbb{C}$ has the density property.

## 5. $\operatorname{Sl}(2, \mathbb{C})$

In this section we will prove Theorem 2.

## Notation and Facts

The complex Lie group $\operatorname{Sl}(2, \mathbb{C})$ will be represented as the set of all $2 \times 2$ matrices with complex entries having determinant 1 . We will write the members of $\mathrm{Sl}(2, \mathbb{C})$ as

$$
A=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \quad \text { with } a d-b c=1
$$

We shall use $a, b, c, d$ as coordinates on $\mathbb{C}^{4}$, in which we will think of $\operatorname{Sl}(2, \mathbb{C})$ as a submanifold. The canonical basis of left invariant vector fields will be employed throughout. These are

$$
\begin{gathered}
X(a, b, c, d)=a \partial_{b}+c \partial_{d}, \quad Y(a, b, c, d)=b \partial_{a}+d \partial_{c} \\
H(a, b, c, d)=a \partial_{a}-b \partial_{b}+c \partial_{c}-d \partial_{d} .
\end{gathered}
$$

The relevant commutation relations are

$$
[H, X]=2 X, \quad[H, Y]=-2 Y, \quad[X, Y]=H
$$

Of course, $X, Y$, and $H$ are $\mathbb{C}$-complete, being left invariant.

Since $X, Y$, and $H$ trivialize the tangent bundle of $\mathrm{Sl}(2, \mathbb{C})$, an arbitrary vector field $V \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$ may be written as

$$
V=V_{X} X+V_{Y} Y+V_{H} H, \quad V_{X}, V_{Y}, V_{H} \in \mathcal{O}(\mathrm{Sl}(2, \mathbb{C}))
$$

We then define

$$
\operatorname{div}(V):=X V_{X}+Y V_{Y}+H V_{H} .
$$

The operator div: $\mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C})) \rightarrow \mathcal{O}(\mathrm{Sl}(2, \mathbb{C}))$ is, up to a constant, the usual divergence operator associated to any left invariant holomorphic 3 -form on $\mathrm{Sl}(2, \mathbb{C})$. Consequently, for any holomorphic function $f$ and vector fields $U$ and $V$, it satisfies:
(i) linearity;
(ii) $\operatorname{div}(f V)=V f+f \operatorname{div} V$; and
(iii) $\operatorname{div}[U, V]=U \operatorname{div} V-V \operatorname{div} U$.

We shall also have occasion to use the right invariant vector fields on $\operatorname{Sl}(2, \mathbb{C})$. The canonical basis is

$$
\begin{gathered}
x=c \partial_{a}+d \partial_{b}, \quad y=a \partial_{c}+b \partial_{d}, \\
h=a \partial_{a}+b \partial_{b}-c \partial_{c}-d \partial_{d} .
\end{gathered}
$$

It is useful to note that

$$
\begin{aligned}
& x=d^{2} X-c^{2} Y+c d H \\
& y=-b^{2} X+a^{2} Y-a b H \\
& h=2 b d X-2 a c Y+(a d+b c) H
\end{aligned}
$$

Finally,

$$
\begin{gathered}
I^{1}(X)=\langle a, c\rangle, \quad I^{1}(Y)=\langle b, d\rangle \\
I^{1}(H)=\left\langle a^{m} b^{k} c^{n} d^{l} \mid m+n-k-l=0\right\rangle .
\end{gathered}
$$

Every orbit of $H$ is biholomorphic to $\mathbb{C}^{*}$, so $I^{2}(H)=I^{1}(H)$. For $X$ and $Y$, the relevant facts about $I^{2}$ are that $X b=a$ and $X d=c$ and that $Y a=b$ and $Y c=$ $d$. We will not need anything about the second integrals of right invariant vector fields, but we will use the facts that $I^{1}(x)=\langle c, d\rangle, I^{1}(y)=\langle a, b\rangle$, and $I^{1}(h)=$ $\left\langle a^{m} b^{k} c^{n} d^{l} \mid m+k-n-l=0\right\rangle$.

## The Volume Density Property

The volume density property for $\mathrm{Sl}(2, \mathbb{C})$ follows immediately from the following theorem.

Theorem 5.1. Every divergence zero polynomial vector field on $\mathrm{Sl}(2, \mathbb{C})$ is divergence zero completely generated.

We shall now prove this theorem. The proof involves many steps, and must be broken up into cases. These cases are isolated according to certain values of an index of monomials. We call this index the $H$-index, and define it as

$$
\operatorname{ind}_{H}\left(a^{m} b^{k} c^{n} d^{l}\right):=m-k+n-l .
$$

Note that $H\left(a^{m} b^{k} c^{n} d^{l}\right)=\operatorname{ind}_{H}\left(a^{m} b^{k} c^{n} d^{l}\right) a^{m} b^{k} c^{n} d^{l}$. A polynomial in $a, b, c$, and $d$ will be called $H$-homogeneous of degree $r$ if the $H$-index of each of its monomials is $r$. We note that $H$-homogeneous polynomials is a concept that descends to $\mathrm{Sl}(2, \mathbb{C})$-that is, when we identify $a d-b c$ and 1 . Let us further point out that, whereas nonzero constants have $H$-index 0 , zero has every integer as its $H$-index. Finally, note that $X$ raises the $H$-index of an $H$-homogeneous polynomial by 2, and that $Y$ lowers the $H$-index of an $H$-homogeneous polynomial by 2:

$$
\begin{aligned}
H\left(X\left(a^{m} b^{k} c^{n} d^{l}\right)\right) & =X H\left(a^{m} b^{k} c^{n} d^{l}\right)+[H, X]\left(a^{m} b^{k} c^{n} d^{l}\right) \\
& =\left(\operatorname{ind}_{H}\left(a^{m} b^{k} c^{n} d^{l}\right)+2\right) X\left(a^{m} b^{k} c^{n} d^{l}\right) \\
H\left(Y\left(a^{m} b^{k} c^{n} d^{l}\right)\right) & =Y H\left(a^{m} b^{k} c^{n} d^{l}\right)+[H, Y]\left(a^{m} b^{k} c^{n} d^{l}\right) \\
& =\left(\operatorname{ind}_{H}\left(a^{m} b^{k} c^{n} d^{l}\right)-2\right) Y\left(a^{m} b^{k} c^{n} d^{l}\right)
\end{aligned}
$$

Finally, we leave it to the reader to check that completeness holds where necessary.
Lemma 5.2. Let $a^{m} b^{k} c^{n} d^{l}$ be a monomial of $H$-index different from -2 . Then there exists a completely generated polynomial vector field of the form

$$
a^{m} b^{k} c^{n} d^{l} X+p(a, b, c, d) H
$$

We shall simultaneously prove the next lemma.
Lemma 5.3. Let $a^{m} b^{k} c^{n} d^{l}$ be a monomial of $H$-index different from 2. Then there exists a completely generated polynomial vector field of the form

$$
a^{m} b^{k} c^{n} d^{l} Y+p(a, b, c, d) H
$$

Proof. We mark the end-of-proof of each case by the symbol $\square$.
Case $1\left(X, \operatorname{ind}_{H} \geq 0\right)$ : Let $m_{1}, m_{2}, n_{1}, n_{2}$ be nonnegative integers such that $m_{1}-k+n_{1}-l=0, m_{1}+m_{2}=m$, and $n_{1}+n_{2}=n$. Then

$$
\left[a^{m_{1}} b^{k} c^{n_{1}} d^{l} H, a^{m_{2}} c^{n_{2}} X\right]=\left(m_{2}+n_{2}+2\right) a^{m} b^{k} c^{n} d^{l} X+p H .
$$

Case $2\left(Y, \operatorname{ind}_{H} \leq 0\right)$ : Let $k_{1}, k_{2}, l_{1}, l_{2}$ be nonnegative integers such that $m-k_{1}+n-l_{1}=0, k_{1}+k_{2}=k$, and $l_{1}+l_{2}=l$. Then

$$
\left[b^{k_{2}} d^{l_{2}} Y, a^{m} b^{k_{1}} c^{n} d^{l_{1}} H\right]=\left(k_{2}+l_{2}+2\right) a^{m} b^{k} c^{n} d^{l} Y+p H
$$

In the remaining cases, the following identities will be very useful:

$$
\begin{align*}
& {\left[a^{m} b^{k_{1}} c^{n} d^{l_{1}} H,\left[b^{k_{2}} d^{l_{2}} Y, a X\right]\right]} \\
& =\left[a^{m} b^{k_{1}} c^{n} d^{l_{1}} H, b^{k_{2}+1} d^{l_{2}} X-a\left(k_{2} a d+l_{2} b c\right) b^{k_{2}-1} d^{l_{2}-1} Y+p H\right] \\
& =\left(1-k_{2}-l_{2}\right)\left(a^{m} b^{k_{1}+k_{2}+1} c^{n} d^{l_{1}+l_{2}} X\right. \\
& \left.\quad-a\left(k_{2} a d+l_{2} b c\right) a^{m} b^{k_{1}+k_{2}-1} c^{n} d^{l_{1}+l_{2}-1} Y\right) \\
& +p H ;  \tag{1}\\
& \begin{aligned}
{\left[a^{m} c^{n} d^{l_{1}} H,\left[d^{l_{2}} Y, c X\right]\right]=} & {\left[a^{m} c^{n} d^{l_{1}} H, d^{l_{2}+1} X-l_{2} c^{2} d^{l_{2}-1} Y+p H\right] } \\
= & \left(1-l_{2}\right)\left(a^{m} c^{n} d^{l_{1}+l_{2}+1} X-l_{2} a^{m} c^{n+2} d^{l_{1}+l_{2}-1} Y\right) \\
& +p H .
\end{aligned}
\end{align*}
$$

Case $3\left(X, \operatorname{ind}_{H} \leq-4, k>0\right)$ : Let $k_{1}, k_{2}, l_{1}, l_{2} \geq 0$ be such that $m-k_{1}+n-l_{1}=0, k=k_{1}+k_{2}+1$, and $l=l_{1}+l_{2}$. Since $m-k+n-l \leq$ $-4,1+m-\left(k_{1}+k_{2}-1\right)+n-\left(l_{1}+l_{2}-1\right) \leq 0$. Thus, using identity (1), we can (via case 2) eliminate the $Y$ component.

Case $4\left(X, \operatorname{ind}_{H} \leq-4, k=0\right)$ : Let $l_{1}, l_{2} \geq 0$ be such that $m+n-l_{1}=0$ and $l=l_{1}+l_{2}+1$. Since $m+n-l \leq-4$, we have $m+(n+2)-\left(l_{1}+l_{2}-1\right) \leq 0$. Thus, using identity (2), we can (again via case 2 ) eliminate the $Y$ component.

Case $5\left(Y, \operatorname{ind}_{H} \geq 4\right)$ : This case can be handled like cases 3 and 4, using appropriate modifications of the identities (1) and (2) and using case 1 instead of case 2. Specifically, one interchanges the roles of $X$ and $Y$, of $a$ and $d$, of $b$ and $c$, of $m$ and $l$, and of $n$ and $k$. The details are left to the interested reader.

Case $6\left(X, \operatorname{ind}_{H}=-1, k>0\right)$ : With $k_{2}=l_{2}=0$, identity (1) takes the form

$$
\left[a^{m} b^{k_{1}} c^{n} d^{l} H,[Y, a X]\right]=a^{m} b^{k_{1}+1} c^{n} d^{l} X+p H
$$

Letting $k=k_{1}+1$ finishes this case.
Case $7\left(X, \operatorname{ind}_{H}=-1, k=0\right)$ : With $l_{2}=0$, identity (2) takes the form

$$
\left[a^{m} b^{k_{1}} c^{n} d^{l} H,[Y, a X]\right]=a^{m} b^{k} c^{n} d^{l_{1}+1} X+p H
$$

Letting $l=l_{1}+1$ finishes this case.
Case $8\left(Y, \operatorname{ind}_{H}=1\right)$ : Again, just use calculations analogous to those of cases 6 and 7.

Case $9\left(X, \operatorname{ind}_{H}=-3\right)$ : Using identities (1) and (2) and case 8, we can eliminate the $Y$ components, which have $H$-index 1 . Notice that, in this case, $1-k_{2}-l_{2} \neq 0$.

Case $10\left(Y, \operatorname{ind}_{H}=3\right)$ : This case is analogous to case 9 .
This completes the proof.
Lemmas 5.2 and 5.3 become false if the index conditions are removed. Fortunately, this is not necessary in order to proceed.

Lemma 5.4. Let $a^{m} b^{k} c^{n} d^{l}$ be an index- -2 monomial. Then there exists a completely generated divergence zero polynomial vector field $V$ of the form

$$
V=a^{m} b^{k} c^{n} d^{l} X+(*) Y+(*) H
$$

Proof. First, let us call a monomial $a^{m} b^{k} c^{n} d^{l}(a, d)$-reduced if either $m$ or $l$ are zero. Every polynomial $p$ on $\mathrm{Sl}(2, \mathbb{C})$ can be written uniquely as a linear combination of $(a, d)$-reduced monomials. Furthermore, $p$ is $(a, d)$-reduced if and only if, for every left invariant vector field $L, L p$ is $(a, d)$-reduced.

Case $1(l>0, m=0)$ : Here $V=b^{k} c^{n} d^{l} X+(*) Y+(*) H$. We thus need only note that

$$
\frac{1}{n+1}\left[b^{k} d^{l-1} Y, c^{n+1} X\right]=b^{k} c^{n} d^{l} X+(*) Y+(*) H
$$

This finishes case 1.
Case $2(l=0)$ : We may assume that $V=a^{m} b^{k} c^{n} X-p Y+(*) H$, where $p$ is an $(a, d)$-reduced, $H$-homogeneous polynomial of $H$-index 2 . Now,

$$
0=\operatorname{div} V=k a^{m+1} b^{k-1} c^{n}-Y p
$$

and so $Y p=k a^{m+1} b^{k-1} c^{n}$. Note that every $(a, d)$-reduced monomial component of $p$ must therefore be of the form $a^{m^{\prime}} b^{k^{\prime}}$. It follows that $n=0$ and that $p$ is a monomial, which must be $C a^{m^{\prime}} b^{k^{\prime}}$. The index conditions then become $m^{\prime}=$ $k^{\prime}+2$ and $m=k-2$. Next, comparing exponents of $Y p$ and $a^{m+1} b^{k-1}$, we see that $m^{\prime}=m+2$. Hence, if $\operatorname{div} V=0$ then $V$ is restricted to be of the form

$$
V=(a b)^{m} b^{2} X-(a b)^{m} a^{2} Y+(*) H
$$

It follows that

$$
V+(a b)^{m} y=(*) H
$$

so $V$ is in fact complete $\bmod H$. This finishes case 2 and thus also the proof of the lemma.

Lemma 5.5. Let $p$ be a nonzero, $H$-homogeneous polynomial of $H$-index 2. Then there is no divergence zero vector field of the form $p Y+q H$.

Proof. Since $Y p$ is of $H$-index 0 , so is $H q$. But since $H$ preserves $H$-index, $q$ is of $H$-index 0 . Hence $H q=0$, so that $Y p=0$. But every nonzero first integral of $Y$ has nonpositive $H$-index. Since $p$ is of $H$-index 2, it must vanish identically.

Proof of Theorem 5.1. Let $V$ be a polynomial vector field of zero divergence. By Lemmas 5.2, 5.3, and 5.4, there is a divergence zero completely generated vector field $W$ such that $V-W=p Y+q H$, where $p$ is an $H$-homogeneous polynomial of $H$-index 2. By Lemma 5.5, $p=0$. Thus $H q=0$, and so $q H$ is complete. We see that $V=W+q H$ is divergence zero completely generated, as desired.

## The Divergence Lemma

Lemma 5.6. Let $V \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$ be a polynomial vector field. Then there exists a completely generated polynomial vector field $W \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$ such that

$$
\operatorname{div} W=\operatorname{div} V
$$

Proof. The image by div of the polynomial vector fields is spanned by the following polynomials:
(i) $\operatorname{div}\left(a^{m} b^{k} c^{n} d^{l} X\right)=a^{m} c^{n} X\left(b^{k} d^{l}\right)$,
(ii) $\operatorname{div}\left(a^{m} b^{k} c^{n} d^{l} Y\right)=b^{k} c^{l} Y\left(a^{m} c^{n}\right)$, and
(iii) $\operatorname{div}\left(a^{m} b^{k} c^{n} d^{l} H\right)=H\left(a^{m} b^{k} c^{n} d^{l}\right)$.

Here $m, k, n$, and $l$ range over all nonnegative integers. We need only show that each of these polynomials is the image by div of a completely generated vector field. To this end, observe that

$$
\operatorname{div}\left[a^{m} c^{n} X, \frac{1}{k+l}\left(k a b^{k-1} d^{l}+l b^{k} c d^{l-1}\right) Y\right]=a^{m} c^{n} X\left(b^{k} d^{l}\right)
$$

and that

$$
\operatorname{div}\left[b^{k} d^{l} Y, \frac{1}{m+n}\left(m a^{m-1} b c^{n}+n a^{m} c^{n-1} d\right) X\right]=b^{k} c^{l} Y\left(a^{m} c^{n}\right)
$$

This takes care of cases (i) and (ii). Case (iii) is only slightly more detailed. To handle it, let $j=m+n-k-l$. If $j=0$, then $H\left(a^{m} b^{k} c^{n} d^{l}\right)=0$ and so there is nothing to do. Suppose that $j>0$. Let $m_{1}, m_{2}, n_{1}, n_{2}$ be nonnegative integers such that:
(a) $m=m_{1}+m_{2}$ and $n=n_{1}+n_{2}$; and
(b) $m_{1}-k+n_{1}-l=0$.

It follows that $m_{2}+n_{2}=j$. Then

$$
\begin{aligned}
& \operatorname{div}\left[a^{m_{1}} b^{k} c^{n_{1}} d^{l} H, \frac{1}{m_{2}+n_{2}}\left(m_{2} a^{m_{2}-1} b c^{n_{2}}+n_{2} a^{m_{2}} c^{n_{2}-1} d\right) X\right] \\
& \quad=a^{m_{1}} b^{k} c^{n_{1}} d^{l} H\left(a^{m_{2}} c^{n_{2}}\right) \\
& \quad=H\left(a^{m} b^{k} c^{n} d^{l}\right)
\end{aligned}
$$

Finally, if $j<0$, let $k_{1}, k_{2}, l_{1}, l_{2}$ be nonnegative integers such that:
(a) $k=k_{1}+k_{2}$ and $l=l_{1}+l_{2}$; and
(b) $m-k_{1}+n-l_{1}=0$.

Then $k_{2}+l_{2}=-j$ and we have

$$
\begin{aligned}
\operatorname{div}[ & \left.a^{m} b^{k_{1}} c^{n} d^{l_{1}} H, \frac{1}{k_{2}+l_{2}}\left(k_{2} a b^{k_{2}-1} d^{l_{2}}+l_{2} b^{k_{2}} c d^{l_{2}-1}\right) Y\right] \\
& =a^{m} b^{k_{1}} c^{n} d^{l_{1}} H\left(b^{k_{2}} d^{l_{2}}\right) \\
& =H\left(a^{m} b^{k} c^{n} d^{l}\right)
\end{aligned}
$$

The reader may confirm directly or via the ideas in Section 3 that all of the vector fields used were complete where required. This completes the proof.

Proof of Theorem 2. Let $U \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$ be a polynomial vector field. By Lemma 5.6 there exists a completely generated vector field $U^{\prime} \in \mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$, which is polynomial, such that $\operatorname{div} U=\operatorname{div} U^{\prime}$. Since $V:=U-U^{\prime}$ is a polynomial vector field with zero divergence, it is (by Theorem 5.1) completely generated. Hence $U=U^{\prime}+V$ is completely generated, and Theorem 2 now follows from the density of polynomial vector fields in $\mathcal{X}_{\mathcal{O}}(\mathrm{Sl}(2, \mathbb{C}))$.

## 6. Elliptic Microspray Manifolds

In this section we explore more fully the density and volume density property on spaces of the form $M \times \mathbb{C}$. The case in which $M$ is a complex Lie group was
already handled in our note [V1]. The proofs of the density theorems in this section are very similar to those in the less general case [V1], and thus will be very sketchy. The main point here is to broaden the class of such complex manifolds $M$ in hopes of giving insight into the density and volume density property.

## Definitions and Examples

Definition 6.1. An elliptic microspray (EM) manifold is a complex manifold $M$ with the property that-for any $V \in \mathcal{X}_{\mathcal{O}}(M)$, compact $K \Subset M$, and $\varepsilon>0$ there exist functions $f_{1}, \ldots, f_{r} \in \mathcal{O}(K)$ and $\mathbb{C}$-completely generated vector fields $X_{1}, \ldots, X_{r}$ satisfying

$$
\left\|V-\sum f_{j} X_{j}\right\|_{K}<\varepsilon
$$

It is also useful to consider slightly more restrictive structures.
Definition 6.2. An $E M V$ ( $V$ for volume) manifold is a pair $(M, \omega)$, where $M$ is a complex manifold and $\omega$ is a holomorphic volume element on $M$, with the property that-for any $V \in \mathcal{X}_{\mathcal{O}}(M)$, compact $K \Subset M$, and $\varepsilon>0$-there exist functions $f_{1}, \ldots, f_{r} \in \mathcal{O}(K)$ and divergence zero completely generated vector fields $X_{1}, \ldots, X_{r}$ satisfying

$$
\left\|V-\sum f_{j} X_{j}\right\|_{K}<\varepsilon
$$

Of course, every EMV manifold is EM. The terminology we have chosen is inspired by that in [G].

Examples. (1) Every complex Lie group $G$ is EMV. Indeed, the left invariant vector fields, which are all complete, parallelize the tangent bundle of $G$, so every vector field can be written in the form $\sum f_{j} V_{j}$, where $f_{j} \in \mathcal{O}(G)$ and $\left\{V_{j}\right\}$ is any fixed basis of $\mathfrak{g}=\operatorname{Lie}(G)$. Moreover, $\operatorname{div} \sum f_{j} V_{j}=\sum V_{j} f_{j}$, so every left invariant vector field has zero divergence.
(2) Every Stein complex homogeneous space is EMV. Indeed, let $G$ be a complex Lie group and $H$ a closed complex subgroup such that $M=H \backslash G=\{H g \mid$ $g \in G\}$ is Stein. The left invariant vector fields on $G$ will project to $M$, as will the left invariant $k$-forms $\left(k=\operatorname{dim}_{\mathbb{C}} M\right)$. Let $V$ be the vector space spanned by the projection to $M$ of the left invariant vector fields on $G$. All of these vector fields have divergence zero with respect to any nonzero volume element coming from a left invariant $k$-form on $G$. Our claim is then proved if we can show that $\mathcal{X}_{\mathcal{O}}(M)=\mathcal{O}(M) \otimes V$. To see the latter, consider the following short exact sequence of coherent sheaves on $M$ :

$$
0 \rightarrow \mathcal{S} \rightarrow \mathcal{O} \otimes V \rightarrow \mathcal{X}_{\mathcal{O}} \rightarrow 0
$$

Here $\mathcal{O}$ is the structure sheaf and $\mathcal{X}_{\mathcal{O}}$ is the tangent bundle sheaf. The sequence gives rise to a long exact sequence in cohomology, a portion of which is

$$
H^{0}(\mathcal{O} \otimes V, M) \rightarrow H^{0}\left(\mathcal{X}_{\mathcal{O}}, M\right) \rightarrow H^{1}(\mathcal{S}, M)
$$

Since $M$ is Stein, $H^{1}(\mathcal{S}, M)=0$ and our claim follows.

## Density Theorems

Our first result is a stable volume density property theorem for EMV manifolds.
Theorem 6.3. If $(M, \omega)$ is an EMV manifold, then $(M \times \mathbb{C}, \omega \wedge d z)$ has the volume density property.

Proof. First, for $X \in \mathcal{X}_{\mathcal{O}}(M)$, let $V=z^{n} X+(*) \partial_{z}$ be a divergence zero vector field. We can assume (by approximation) that $X=\sum \varphi_{j} Y_{j}$, with $Y_{j} \in \mathcal{X}_{\mathcal{O}}(M, \omega)$ divergence zero $\mathbb{C}$-completely generated. Now

$$
\left[1 /(n+1) z^{n+1} Y_{j}, \varphi_{j} \partial_{z}\right]=z^{n} \varphi_{j} Y_{j}+(*) \partial_{z}
$$

is clearly divergence zero $\mathbb{C}$-completely generated; hence $V$ is divergence zero $\mathbb{C}$-completely generated modulo $\partial_{z}$. That is, there exists a holomorphic vector field $W$ that is divergence zero $\mathbb{C}$-completely generated and has the property that $V-W=\psi(x, z) \partial_{z}$. But then $0=\operatorname{div}(V-W)=\partial_{z} \psi$, so that $V-W$ is complete. Hence $V=W+(V-W)$ is divergence zero $\mathbb{C}$-completely generated. Since every divergence zero vector field can be approximated by sums of vector fields of the same form as $V$, we are done.

The next result is that EM manifolds with holomorphic volume elements are stably EMV.

Proposition 6.4. If $M$ is an EM manifold and $\omega$ is a nonvanishing holomorphic volume element on $M$, then $(M \times \mathbb{C}, \omega \wedge d z)$ is EMV.

We shall need the following lemma.
Lemma 6.5. Let $M$ and $\omega$ be as in Proposition 6.4. If $X \in \mathcal{X}_{\mathcal{O}}(M)$ is ( $\left.\mathbb{C}-\right)$ completely generated, then there exists $\tilde{X} \in \mathcal{X}_{\mathcal{O}}(M \times \mathbb{C}, \omega \wedge d z)$, which is divergence zero completely generated, such that

$$
\tilde{X}-X=(*) \partial_{z} .
$$

Proof. First note that, if $X \in \mathcal{X}_{\mathcal{O}}(M)$ is complete, then so is $X-z\left(\operatorname{div}_{\omega} X\right) \partial_{z}$. Moreover, the latter has zero $(\omega \wedge d z)$-divergence. Next notice that

$$
X+(*) \partial_{z}+Y+(*) \partial_{z}=X+Y+(*) \partial_{z}
$$

and that

$$
\left[X+(*) \partial_{z}, Y+(*) \partial_{z}\right]=[X, Y]+(*) \partial_{z} .
$$

The lemma follows easily from these facts.
Proof of Proposition 6.4. Let $X \in \mathcal{X}_{\mathcal{O}}(M \times \mathbb{C}, \omega \wedge d z)$ be written as $X=$ $\sum z^{j} V_{j}+(*) \partial_{z}$, where $V_{j} \in \mathcal{X}_{\mathcal{O}}(M)$. By approximation, we may assume that the sum is finite. Since $M$ is EM, we may write (again, up to approximation) $V_{j}=\sum_{k} f_{j k} S_{j k}$, where the $S_{j k} \in \mathcal{X}_{\mathcal{O}}(M)$ are $\mathbb{C}$-completely generated. Now, for each $S_{j k}$, the lemma guarantees a divergence zero completely generated $\tilde{S}_{j k} \in$ $\mathcal{X}_{\mathcal{O}}(M \times \mathbb{C}, \omega \wedge d z)$ such that $S_{j k}-\tilde{S}_{j k}=(*) \partial_{z}$. It follows that (up to approximation)

$$
X=\sum_{j k} z^{j} f_{j k} \tilde{S}_{j k}+(*) \partial_{z}
$$

which is exactly what was needed.
Using main result I. 3 in [V1], one immediately obtains the following.
Corollary 6.6. If $M$ is a Stein EMV space, then $M \times \mathbb{C}$ has the density property. If $M$ is a Stein EM space and $M$ admits a holomorphic volume element, then $M \times \mathbb{C}^{2}$ has the density property.

## 7. A Question

The results in Section 6 suggest the following natural question:
Is there a difference between the volume density property and EMV?
To date, in all the examples for which we have been able to settle this question, the answer is No. If this answer can be established in general, it would represent a major breakthrough. However, it is by no means clear what the answer is. Again, one needs candidates for testing. We propose one now. Let

$$
\Sigma^{3}:=\left\{(a, b, c, d) \in \mathbb{C}^{4} \mid a^{2} d-b c=1\right\}
$$

$\Sigma^{3}$ is a smooth subvariety of $\mathbb{C}^{4}$ and is also a branched double cover of $\operatorname{Sl}(2, \mathbb{C})$. Moreover, $\Sigma^{3}$ admits some interesting complete vector fields:

$$
\begin{gathered}
X=a^{2} \partial_{b}+c \partial_{d}, \quad Y=b \partial_{a}+2 a d \partial_{c} \\
H=a \partial_{a}-2 b \partial_{b}+2 c \partial_{c}-2 d \partial_{d}
\end{gathered}
$$

correspond to the left invariant vector fields of $\operatorname{Sl}(2, \mathbb{C})$, and

$$
\begin{gathered}
\xi=a^{2} \partial_{c}+b \partial_{d}, \quad \eta=c \partial_{a}+2 a d \partial_{b} \\
\theta=a \partial_{a}+2 b \partial_{b}-2 c \partial_{c}-2 d \partial_{d}
\end{gathered}
$$

correspond to the right invariant vector fields of $\operatorname{Sl}(2, \mathbb{C})$. Since $\Sigma^{3}$ is 3-dimensional, we expect some relations between the left and right vector fields. A calculation shows that

$$
\begin{gathered}
\xi=-b^{2} X+\frac{1}{2} a^{3} Y-\frac{1}{2} a^{2} b H, \quad \eta=2 a d^{2} X-c^{2} Y+a c d H \\
\theta=4 b d X-2 a c Y+\left(a^{2} d+b c\right) H
\end{gathered}
$$

We define a volume element $\Omega$ on $\Sigma^{3}$ as follows. Set

$$
\begin{gathered}
\Omega_{X}=d \delta_{b}-b \delta_{d}, \quad \Omega_{Y}=\frac{1}{2}\left(a \delta_{c}-2 c \delta_{a}\right) \\
\Omega_{H}=\frac{1}{2}\left(2 a d \delta_{a}+c \delta_{b}-b \delta_{c}-a^{2} \delta_{d}\right)
\end{gathered}
$$

and define $\Omega=\Omega_{X} \wedge \Omega_{Y} \wedge \Omega_{H}$. Here, $\delta_{a}\left(\partial_{x}\right)=0$ if $x=b, c, d$ and 1 if $x=a$, and similarly for $\delta_{b}, \delta_{c}$, and $\delta_{d}$. One can easily compute the following:

$$
\begin{gathered}
{[H, X]=4 X, \quad[H, Y]=-3 Y, \quad[X, Y]=a H} \\
\operatorname{div}(X)=\operatorname{div}(Y)=0
\end{gathered}
$$

It follows that $\operatorname{div}(a H)=0$ and hence that $\operatorname{div}(H)=1$. Thus, since $a H$ vanishes when $a=0$, we need more than just $X, Y$, and $H$ to prove that $\Sigma^{3}$ is EMV. However, this is indeed the case.

Proposition 7.1. $\quad \Sigma^{3}$ is EMV.
Proof. It suffices to show that $H$ can be written as a sum $\sum f_{j} V_{j}$ with the $V_{j}$ generated by $X$ and $Y$. To this end,

$$
\begin{aligned}
a d[X, Y]+c[Y,[Y, X]]-3 a c Y & =a^{2} d H-c[Y, a H]-3 a c Y \\
& =a^{2} d H-c((Y a) H-a[Y, H])-3 a c Y \\
& =\left(a^{2} d-b c\right) H+3 a c Y-3 a c Y \\
& =H
\end{aligned}
$$

Moreover, we have been able to prove (with considerable difficulty) that if $\Sigma^{3}$ has the volume density property, then it has the density property. Nevertheless, the combinatorics arising in attempts to prove the volume density property by the methods of Sections 4 and 5 become too cumbersome for us to handle.

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