Integration of analytic differential systems with singularities and some applications to real submanifolds of C^n *

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1. Introduction.

A module *D* of analytic vector fields on \mathbb{R}^n defines at each $y \in \mathbb{R}^n$ a subspace $D(y) = \{X(y) : X \in D\}$ of the tangent space to \mathbb{R}^n at y. A real-analytic submanifold *M* of \mathbb{R}^n is an *integral manifold of D* if

$$T_y M = D(y)$$
 for all $y \in M$, (1.1)

where T_yM is the tangent space to M at y. In [6] Nagano proved that if D is closed under the Lie bracket then through each point there passes a unique integral manifold of D. This result extends the classical Frobenius theorem, which assumes in addition that D(y) has constant dimension. In dropping this hypothesis, Nagano relies on the analyticity. The classical theorem also holds in the C^{∞} category and in [6] Nagano gives a simple C^{∞} counterexample to his result.

This paper contains 1) a new proof of Nagano's theorem in a formulation which describes the integral manifold directly in terms of D, 2) a sharpened form of the theorem giving necessary and sufficient conditions at p for the existence of an integral manifold through p, and 3) some applications of these results to the local geometry of real-analytic submanifolds of a complex manifold. In particular, it is shown that a point p on a real-analytic CR submanifold M is not of finite weight [1] if and only if there is a complex submanifold of M of maximum dimension through p.

The proofs given here are almost entirely algebraic and make no use of differential equations. They appear to be new even in the classical case where D(y) has constant dimension. Besides a simple and standard majorization argument and advanced calculus, one needs only the standard Weierstrass division theorem [3, Satz 1, p. 23]. All definitions are within the real-analytic category,

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and unless otherwise stated are to be interpreted as germs at a fixed point p of \mathbb{R}^n (recall that there is a standard interpretation of (1.1) for germs). The ring of (germs at p of) real-analytic functions is denoted \mathcal{C}^{ω} and the \mathcal{C}^{ω} -module of real-analytic vector fields is denoted $T\mathbb{R}^n$.

2. Integration Theorems.

A submanifold M defines the ideal I_M of all functions in \mathcal{C}^{ω} which vanish on M, and $M \supset N$ if and only if $I_M \subset I_N$. If I is an ideal in \mathcal{C}^{ω} then the variety V(I) is the set of points where all functions in I vanish. The ideal generated by functions f_{l+1}, \dots, f_n is denoted (f_{l+1}, \dots, f_n) . If these functions are independent at p then $M = V(f_{l+1}, \dots, f_n)$ if and only if

$$\boldsymbol{I}_{\boldsymbol{M}} = (f_{l+1}, \cdots, f_n). \tag{2.1}$$

Of course, in this case M has dimension l. A vector field $X \in T\mathbb{R}^n$ is tangent to M if $X(y) \in T_y M$ for all $y \in M$. It is elementary that this is equivalent to $Xf_j \in I_M$ for $j = l+1, \dots, n$. In view of (2.1) and the Leibniz rule X(fg)= fXg + gXf,

X is tangent to M if and only if
$$I_M$$
 is closed under X (2.2)

(which, of course, means that $Xf \in I_M$ for every $f \in I_M$). Since \mathcal{C}^{ω} is Noetherian D is finitely generated. Choose now a set X_1, \dots, X_k of generators for D, and for each k-index $\alpha = (\alpha_1, \dots, \alpha_k)$ of nonnegative integers let $X^{\alpha} = X_1^{\alpha_1} \cdots X_k^{\alpha_k}$. In X^{α} the X_j 's always occur in serial order, and as usual X^0 is defined by $X^0f = f$.

If M is an integral manifold of D, (1.1) and (2.2) show that I_M is invariant under D. Therefore, I_M is contained in

$$I_{D} = \{g \in \mathcal{C}^{\omega} : X^{\alpha}g(0) = 0 \quad \text{for all} \quad \alpha\}.$$
(2.3)

It follows from the Leibniz rule that I_D is an ideal in C^{ω} . If D is bracketclosed, the Leibniz rule and Lemma 2.3 below show that I_D is independent of a choice of generators for D. This ideal provides an invariant means for utilizing the basic principle that an analytic function is determined by its derivatives at a fixed point. The following is equivalent to Nagano's theorem 1 in [6].

THEOREM 2.1. If D is a bracket-closed submodule of $T\mathbf{R}^n$ then $V(\mathbf{I}_D)$ is the unique integral manifold of D through p.

PROOF. In order that $V(I_D)$ is a manifold, there must exist independent functions f_{l+1}, \dots, f_n such that

$$\boldsymbol{I}_{D} = (f_{l+1}, \cdots, f_{n}), \qquad (2.4)$$

and if $V(I_D)$ is an integral manifold then $l=\dim D(p)$ must hold. It will be shown first that there exist such functions, and later that $V(I_D)$ is indeed an integral manifold.

It may be assumed that p=0. In terms of analytic coordinates $x=(x_1, \dots, x_n)$ each $X_i = \sum_{j=1}^n a_{ij}\partial/\partial x_j$, where the matrix (a_{ij}) has rank $l=\dim D(0)$ at 0. By a reindexing, it may be assumed that $(a_{ij})_{i,j=1}^l$ is non-singular at 0. Since elementary row operations over \mathcal{C}^{ω} can be achieved by linear combination of the X_i , it may be assumed that the generators X_i of D have the form

$$X_{i} = \partial/\partial x_{i} + \sum_{j=l+1}^{n} a_{ij} \partial/\partial x_{j}, \quad i=1, \dots, l, \text{ and}$$

$$X_{i} = \sum_{j=l+1}^{n} a_{ij} \partial/\partial x_{j}, \quad i=l+1, \dots, k,$$
(2.5)

and $a_{ij}(0)=0$, i, j>l. This has the usual geometric interpretation that an integral manifold is spread over the x_1, \dots, x_l coordinate space. Since it must be the graph of a function of these variables it is natural to seek defining functions of the form

$$f_i = x_i - g_i(x_1, \dots, x_l), \quad i = l+1, \dots, n,$$
 (2.6)

where each g_i is analytic.

An analytic function depending only on x_1, \dots, x_l will be written $g = \sum_{\alpha}' a_{\alpha} x^{\alpha}$, where the prime means that summation extends only over those α such that $\alpha_{l+1} = \dots = \alpha_n = 0$. By (2.5) such a function satisfies $X_i g = \partial g / \partial x_i$, $i \leq l$, so that

$$X^{\alpha}g(0) = \partial^{|\alpha|}g/\partial x^{\alpha}(0) = a_{\alpha}\alpha$$
 (2.7)

if $\alpha_{l+1} = \cdots = \alpha_n = 0$.

Therefore, in order that (2.6) defines functions in I_D it is necessary that

$$g_i = \sum'(\alpha !)^{-1} X^{\alpha}(x_i)(0) x^{\alpha} .$$
(2.8)

LEMMA 2.2. The series (2.8) converges.

This will be proved later.

In the classical Frobenius theorem D(y) has constant dimension, which is equivalent to k=l. In this case membership of f_i in I_D is immediate from (2.6), (2.7), and (2.8). When k>l this fact also depends on the property that

for any f,
$$X_i f(0) = 0$$
 if $i > l$, (2.9)

and on Lemma 2.3 below. A *polynomial* is any linear combination over C^{ω} of operators of the form X^{α} as defined above. If α and β are multiindices then $\alpha + \beta$ is their usual vector sum and $|\alpha| = \alpha_1 + \cdots + \alpha_n$. For $1 \leq j \leq l$, $e^j = (0, \cdots, 0, 1, 0, \cdots, 0)$ (the 1 in the *j*-th place), and $\alpha \pm e^j$ is written $\alpha \pm j$.

LEMMA 2.3. If D is bracket closed then for any α , $\beta X^{\alpha}X^{\beta}-X^{\alpha+\beta}$ is a polynomial of degree less than $|\alpha+\beta|$.

This is also proved later. It is now used to show by induction on $|\gamma|$ that

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$$X^{r}f_{j}(0)=0, \quad j=l+1, \cdots, n.$$
 (2.10)

Relations (2.7) and (2.9) imply the case $|\gamma|=1$. Now assume m>1 and (2.10) holds for all $|\gamma| < m$, and suppose $|\gamma|=m$. If $\gamma_{l+1}=\cdots=\gamma_n=0$ then (2.10) holds as shown above. Otherwise $\gamma=\beta+i$ for some i>l and Lemma 2.3 with $\alpha=e^i$ shows that $X^{\gamma}=X_iX^{\beta}+(\text{polynomial of degree}<m)$. Thus (2.10) follows for γ by the induction hypothesis and (2.9). This proves that each $f_i \in I_D$ and hence

$$(f_{l+1}, \cdots, f_n) \subset I_D.$$

$$(2.11)$$

To show (2.4), suppose that f is any function in C^{ω} . By the Weierstrass Division Theorem [3] $f=a_{l+1}f_{l+1}+r_{l+1}$ where r_{l+1} is a polynomial in x_{l+1} of degree less than the order of the zero $x_{l+1} \rightarrow f_{l+1}(0, \dots, 0, x_{l+1}, 0, \dots, 0)=x_{l+1}$ has at 0. In other words, r_{l+1} is independent of x_{l+1} . So is f_{l+2} , and hence r_{l+1} may be divided by f_{l+2} in the subring of functions independent of x_{l+1} to obtain $f=a_{l+1}f_{l+1}+a_{l+2}f_{l+2}+r_{l+2}$, where r_{l+2} is independent of x_{l+1} and x_{l+2} . Continuing in this way will achieve

$$f = a_{l+1}f_{l+1} + \dots + a_n f_n + r_n \tag{2.12}$$

where r_n depends only on x_1, \dots, x_l . If also $f \in I_D$, then so does r_n by (2.11). By (2.7), $r_n=0$. This proves (2.4), and hence $M=V(I_D)=V(f_{l+1}, \dots, f_n)$ is a manifold.

Next, it is claimed that

$$I_D$$
 is closed under D . (2.13)

This is proved by an easy induction on Lemma 2.3, this time with $|\beta|=1$. It follows that M is an integral manifold, for (2.2) and (2.13) imply that each vector field in D is tangent to M. Hence $D(y) \subset T_y M$ for all y on M near 0. This implies (1.1) since dim $D(y) \ge \dim D(0) = l = \dim T_y M$ for all y near 0.

It remains to prove the uniqueness. If N is any integral manifold through 0, (1.1) implies that $D \subset TN$, the module of vector fields tangent to N. Therefore $I_M = I_D \supset I_{TN} \supset I_N$, in which the containments are obvious from the definition of I_{TN} as in (2.3) (with TN replacing D and any choice of generators for TN). It follows that $M \subset N$, and since both manifolds have dimension l they are equal. This completes the proof of Theorem 2.1.

PROOF OF LEMMA 2.2. This simple majorization argument is given in detail (and greater generality) in [4], so it is only sketched here. There exist $\rho > 0$ and M > 0 such that the geometric series $G(t) = M(1-(t/\rho))^{-1}$ majorizes f and all a_{ij} at 0, and such that $X_j(x_i)$ is majorized by ΔG at 0, where Δ is the operator $\Delta = G(t) d/dt$. It follows inductively that each $X^{\alpha}(x_i)$ is majorized by $\Delta^{|\alpha|}G =$ $M^{|\alpha|+1}\rho^{-|\alpha|}(1\cdot3\cdot5\cdots(2|\alpha|-1))(1-(t/\rho))^{2|\alpha|+1}$. Therefore $|(X^{\alpha}(x_i))(0)| \leq M^{|\alpha|+1}\rho^{-|\alpha|}$ $(1\cdot3\cdot5\cdots(2|\alpha|-1))$, and the ratio test now gives convergence of (2.8). PROOF OF LEMMA 2.3. If $|\alpha + \beta| = 2$ the lemma is either trivial ($\alpha = 0$ or $\beta = 0$), or is just a restatement of the bracket closure hypothesis, which is equivalent to

$$X_{i}X_{j} - X_{j}X_{i} = \sum_{q=1}^{k} a_{ij}^{q}X_{q}.$$
 (2.14)

Assume $m \ge 2$ and the lemma holds for all α , β with $|\alpha + \beta| \le m$. It follows that if $|\alpha| + q \le m$ and P is a polynomial of degree $\le q$ then $X^{\alpha}P$ and PX^{α} are polynomials of degree $\le |\alpha| + q$. This is obvious for PX^{α} , and for $X^{\alpha}P$ a simple induction on the Leibniz rule will account for the differentiations of the coefficients of P by X^{α} .

Now suppose that $|\alpha+\beta|=m+1$. It may be assumed that $\beta\neq 0$.

The case $|\beta|=1$: Then $\beta=e^{j}$. Let *i* be the largest integer *r* such that $\alpha_{r}\neq 0$. If $i\leq j$ then $X^{\alpha}X^{\beta}=X^{\alpha}X_{j}=X^{\alpha+j}=X^{\alpha+\beta}$. Otherwise, by (2.14), $X^{\alpha}X^{j}=X^{\alpha-i}X_{i}X_{j}=X^{\alpha-i}X_{j}X_{i}+X^{\alpha-i}P$ where *P* is a polynomial of degree 1. Since $|\alpha-i|=m-1$ the inductive hypothesis yields $X^{\alpha-i}X_{j}=X^{\alpha-i+j}+Q$, where *Q* is a polynomial of degree at most m-1. Thus $X^{\alpha}X_{j}=X^{\alpha-i+j}X_{i}+QX_{i}+X^{\alpha-i}P=X^{\alpha+j}+(\text{polynomial of degree }\leq m)$ because i>j and the inductive assumption applies to $QX_{i}+X^{\alpha-i}P$. This establishes the case $|\beta|=1$.

The general case $|\beta| \ge 1$: Then $\beta = \gamma + e^j$, where j is the largest integer r such that $\beta_{\tau} \ne 0$, and $|\alpha + \gamma| = m$. Hence $X^{\alpha} X^{\beta} = X^{\alpha} X^{\gamma} X_j = X^{\alpha + \gamma} X_j + P X_j$, where P is a polynomial of degree at most m-1. By Case 1 applied to $X^{\alpha + \gamma} X_j$ and the inductive assumption applied to PX_j , $X^{\alpha} X^{\beta} = X^{\alpha + \gamma + j} + Q = X^{\alpha + \beta} + Q$, where Q is a polynomial of degree at most m. This proves Lemma 2.3.

The condition (1.1) can be replaced by

$$TM = D + O(M) \tag{2.15}$$

where O(M) denotes all vector fields vanishing on M and TM is the module of vector fields tangent to M. Condition (2.15) appears stronger than (1.1) (after all $\operatorname{sp}_{R}\{x(d/dx)\} = \operatorname{sp}_{R}\{x^{2}d/dx\}$ for each x in R but xd/dx is not in the C^{ω} -module spanned by $x^{2}(d/dx)$), but in fact it is equivalent. For certainly (2.15) implies (1.1) and if (1.1) holds then clearly $D \subset TM$. If $X \in TM$ there exist by (2.5) functions c_{1}, \dots, c_{l} in C^{ω} such that $Y = X - \sum_{i=1}^{l} c_{i}X_{i}$ is free from terms involving $\partial/\partial x_{1}, \dots, \partial/\partial x_{l}$. Since Y is also tangent to M, it is easy to see that $Y \in O(M)$, and (2.15) is proved.

If p is fixed in Theorem 2.1, the bracket-closure is not necessary. For example, let $X=\partial/\partial x$ and $Y=\partial/\partial y+xz\partial/\partial z$ in \mathbb{R}^3 with coordinates (x, y, z). Then V(z) is an integral manifold through O of $D=\sup\{X, Y\}$, but $[X, Y]=z\partial/\partial z \in D$. Let \tilde{D} denote the smallest bracket-closed submodule containing a given submodule D.

COROLLARY 2.4. There exists an integral manifold of D through p if and

only if

$$\widetilde{D}(p) = D(p), \qquad (2.16)$$

and an integral manifold through p is unique.

PROOF. If M is any manifold the module TM is bracket-closed. Thus if M is an integral manifold then $D \subset TM$ implies

$$D \subset \widetilde{D} \subset TM = D + O(M) \tag{2.17}$$

which in turn implies (2.16). It also shows that M is an integral manifold of \tilde{D} , which implies its uniqueness by that part of Theorem 2.1.

For the converse, apply Theorem 2.1 to \tilde{D} and note that its integral manifold M through p satisfies $T_y M = \tilde{D}(y) \supset D(y)$ for all y on M near p. Because of (2.16) $T_p M = D(p)$. Hence for y on M near p, dim $D(y) \ge \dim T_y M$, and consequently $T_y M = D(y)$. This proves the Corollary.

In the example above [X, [X, Y]] = 0 = [Y, [X, Y]] so $\tilde{D} = \operatorname{sp} \{X, Y, [X, Y]\}$ and $\tilde{D}(0) = \operatorname{sp}_{R} \{\partial/\partial x, \partial/\partial y\} = D(0)$.

3. Existence of complex submanifolds of a real-analytic submanifold.

Let $C_{\mathcal{C}}^{\omega} = C^{\omega} + iC^{\omega}$ be the (Noetherian) ring of complex-valued real-analytic functions and $CTR^n = TR^n + iTR^n$ the complexified space of vector fields. If Mis a real analytic submanifold then I_M now denotes the ideal of all functions in $C_{\mathcal{C}}^{\omega}$ vanishing on M, and CTM = TM + iTM the complexified space of vector fields tangent to M. It is clear that CTM is a bracket-closed module over $C_{\mathcal{C}}^{\omega}$ and that it is defined by (2.2).

There is a natural conjugation $Z=X+iY \rightarrow \overline{Z}=X-iY$ on CTR^n , where X and Y are in TR^n , and a submodule D of CTR^n is real if $D=\overline{D}$. Defining $\operatorname{Re} D=D\cap TR^n$ it is easy to see that D is real if and only if $D=\operatorname{Re} D+i\operatorname{Re} D$, and that a real D is bracket closed if and only if $\operatorname{Re} D$ is bracket closed. Let \widetilde{D} be the bracket closure of D in CTR^n , and $\operatorname{Re} D$ the bracket closure of $\operatorname{Re} D$ in TR^n (as in Section 2). If D is real then it is easily seen that $\widetilde{D}=\operatorname{Re} D+i\operatorname{Re} D$, so that \widetilde{D} is also real and $\operatorname{Re} \widetilde{D}=\operatorname{Re} D$.

A submanifold M is an integral manifold of D if

$$CTM = D + O(M), \qquad (3.1)$$

where O(M) is the C_c^{ω} module of vector fields vanishing on M. Since CTM and O(M) are real, so is D if (3.1) holds. It is easy to see that (3.1) holds if and only if $TM = \operatorname{Re} D + \operatorname{Re} O(M)$, so that M is an integral manifold of D if and only if it is for $\operatorname{Re} D$.

Therefore the extension of Theorem 2.1 to submodules of CTR^n is the following.

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THEOREM 3.1. If a real submodule D of CTR^n is bracket closed then it has a unique integral manifold through p.

Moreover, by the above remarks, Corollary 2.4 is true verbatim for real submodules of CTR^n , where \tilde{D} is interpreted as above.

If \mathbb{R}^{2n} is the real space underlying \mathbb{C}^n , one writes $CT\mathbb{C}^n$ for $CT\mathbb{R}^{2n}$. As usual, it carries another complex structure induced from \mathbb{C}^n , whose action is denoted by J. The module of *complex tangent* vector fields to a submanifold M is $H = \mathbb{C}TM \cap J\mathbb{C}TM$, the largest J-invariant submodule of $\mathbb{C}TM$. The submanifold M is complex if and only if $H = \mathbb{C}TM$; i.e. $\mathbb{C}TM$ is J-invariant. Since O(M) is J-invariant, it follows from (3.1) that

It is shown in [2] that if M is any submanifold then $N_2 = \{X \in H : [X, H] \subset H\}$ is real, *J*-invariant, and bracket closed $(N_2$ can often be viewed as the nullspace of the E.E. Levi form of M). It is also shown in [2] that there exists a chain (which in general is strictly decreasing) $N_2 \supset N_3 \supset \cdots \supset S_0 = \bigcap_{j=2}^{\infty} N_j$ of real, *J*-invariant, bracket closed submodules. The following is thus a consequence of Theorem 3.1.

THEOREM 3.2. A real analytic submanifold M has a chain of successively finer partitions by the complex integral manifolds of the submodules N_i .

It is possible for M to have a complex submanifold larger than the integral manifold of N_2 . In [7], Sommer showed that $M = \{|z_1|^2 + |z_2|^2 - |z_3|^2 - 1 = 0\}$ has 1-dimensional complex submanifolds but $N_2 = 0$. It is therefore of interest to seek an integral manifold of H through a fixed point p on M. In case M is CR, this would be the largest complex submanifold of M through p. The answer is immediate from Corollary 2.4 and (3.2).

THEOREM 3.3. There exists a (unique) integral manifold of H through $p \in M$ if and only if $H(p) = \tilde{H}(p)$. This integral manifold is a complex submanifold of M.

All statements of Theorem 3.3 but the last are proved above. If S is an integral manifold of H through p then $H \subset CTM$ implies $I_S = I_H \supset I_{CTM} = I_M$, so that $S \subset M$. It is easy to see that \tilde{H} is the module spanned over C_c^{∞} by the set H' of all repeated commutators of elements of H. Therefore $\tilde{H}(p) = H(p)$ if and only if H'(p) = H(p). This is a kind of infinitely high-order flatness condition at p. It means that p is not a point of finite type in the sense of Bloom and Graham [1] or Kohn [5], and it causes behavior in contrast with the case of finite type. Thus if M is a real analytic hypersurface satisfying Theorem 3.3 there exists a neighborhood U of p such that the integral manifold of H through z is the zero set of a holomorphic function f. Then 1/f is a holomorphic function on U-M which cannot be continued across p. The local Levi problem [5, p. 528] is thus solvable near p in this case. However, such a point is clearly

never a local peak point [5, p. 540].

References

- T. Bloom and I. Graham, On "type" conditions for generic real submanifolds of Cⁿ, Invent. Math., 40 (1977), 217-243.
- [2] M. Freeman, Local biholomorphic straightening of real submanifolds, Ann. of Math., 106 (1977), 319-352.
- [3] H. Grauert and R. Remmert, Analytische Stellenalgebren, Springer-Verlag, Berlin, 1971.
- [4] W. Gröbner and H. Knapp, Contributions to the method of Lie Series, Bibliographisches Institut, Mannheim, 1967.
- [5] J.J. Kohn, Boundary behavior of \$\overline{\delta}\$ on weakly pseudo-convex manifolds of dimension two, J. Differential Geometry, 6 (1972), 523-542.
- [6] T. Nagano, Linear differential systems with singularities and an application to transitive Lie algebras, J. Math. Soc. Japan, 18 (1966), 398-404.
- [7] F. Sommer, Komplex-analytische Blätterung reeller Hyperflächen im Cⁿ, Math. Ann., 137 (1959), 392-411.

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