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### On completeness of invariant measures defined by differential equations

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## On completeness of invariant measures defined by differential equations;

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#### I. — Introduction.

Let the real, first order, non-singular, ordinary differential system

(1) 
$$\frac{dx}{dt} = f(x, y), \quad \frac{dy}{dt} = g(x, y) \quad (f^2 + g^2 > 0),$$

with

$$f, g \in C^{(\alpha)}$$
  $(\alpha = 1, 2, 3, ..., \infty, A)$  (1)

be defined in an open plane set R. Then through each point  $P_0$ :  $(x_0, y_0) \in \mathbb{R}$  there exists an unique solution curve

$$x = x(t; x_0, y_0), \quad y = y(t; x_0, y_0),$$

initiating at  $P_0$  for t = 0, and defined for some maximal « time » interval

$$\tau_{-}(\mathbf{P}_0) < t < \tau_{+}(\mathbf{P}_0).$$

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<sup>(1)</sup>  $f(x_1, x_2, \ldots, x_n) \in C^{(0)}$  in R, an open set of Euclidean *n*-space, in case f(x) is continuous, in *n* variables in R :  $f(x_1, x_2, \ldots, x_n) = f(x) \in C^{(K)}$ (K = 1, 2, ...) in R in case all the partial derivatives of f(x) up to and including those of order K exist (are finite) and are continuous in R.  $f(x) \in C^{(\infty)}$  in case all partial derivatives of f(x) exist and are continuous in R.  $f(x) \in C^{(n)}$  in case f(x) is analytic in R, that is, near each point in R, f(x) has an absolutely convergent, real, *n*-variable power series representation. We define  $o < 1 < 2 < ... < \infty < A$  and also  $\infty \pm K = \infty$  and  $A \pm K = A$ .

The subset of Euclidean 3-space in which  $x(t; x_0, y_0)$  and  $y(t; x_0, y_0)$ are defined is open and therein these functions, as well as  $\frac{\partial x}{\partial t}$  and  $\frac{\partial y}{\partial t}$ , are in class  $C^{(\alpha)}$ . The transformations  $T_t$ 

$$(x_0, y_0) \rightarrow x(t; x_0, y_0), y(t; x_0, y_0),$$

from an open subset  $R_0 \subset R$  onto open subsets  $R_i \subset R$ , from a one-parameter local group, or stream, of  $C^{(\alpha)}$ -homeomorphisms (<sup>2</sup>).

Suppose the ordinary differential system (1) were exact, that is,  $\frac{\partial f}{\partial x} + \frac{\partial g}{\partial y} \equiv 0$  and R were a simply-connected region (open, connected set (3)). Then  $\{T_i\}$  forms a measure true stream, that is,  $T_i$ is measurability preserving ( $T_i$  and  $T^{-1}$  preserve Lebesgue measurability of sets) and also  $T_i$  preserves the magnitude of the measure. In this case, there exists a stream function  $\psi(x, y) \in C^{(\alpha+1)}$  in R such that  $f = \frac{\partial \psi}{\partial y}$ ,  $g = -\frac{\partial \psi}{\partial x}$  and  $\psi(x, y)$  is a principal integral (4) of the first order, partial differential equation

(2) 
$$f(x, y)\frac{\partial \psi}{\partial x} + g(x, y)\frac{\partial \psi}{\partial y} = 0.$$

That is,  $\psi(x, \gamma)$  is a solution of (2) which is constant on no open set and therefore the class of all solutions of (2) are precisely those functions in C<sup>(1)</sup> in R which are functionally dependent on  $\psi(x, \gamma)$ .

Even if the ordinary differential system (1) is not exact, there may exist a principal integral  $\psi(x, \gamma) \in C^{(2)}$  of the corresponding partial differential equation (2). Then

$$\mu(x, y) = \left[\frac{\psi_x^2 + \psi_y^2}{f^2 + g^2}\right]^{\frac{1}{2}} \in \mathbf{C}^{(1)}$$

<sup>(&</sup>lt;sup>2</sup>) A homeomorphism T of an open set  $R_1$  of Euclidean *n*-space onto a second open set  $R_2$  of the same space is called a  $C^{(\beta)}$ -homeomorphism,  $\beta \equiv 0, 1, 2, ..., \infty$ , A in case both T and  $T^{-1}$  are expressed by functions of class  $C^{(\beta)}$ .

<sup>(3)</sup> For notation, see E. HOPF, Ergodentheorie (Berlin, 1937).

<sup>(\*)</sup> E. KAMKE, Differentialgleichungen Reeller Funktionen (Chelsea, 1947), p. 323.

is a non-negative integrating factor for (1), that is,

$$\mu f = \psi_y, \qquad \mu g = -\psi_x.$$

In case (1) is exact then  $\mu(x, y)$  is a constant.

Because of the relation (5).

(3) 
$$\frac{d}{dt}\left[\iint_{\mathbf{R}_0}\mu(x, y)\frac{\partial(x, y)}{\partial(x_0, y_0)}dx_0\,dy_0\right] = \iint_{\mathbf{R}_0}\left[(\mu f)_x + (\mu g)_y\right]dx_0\,dy_0 = 0$$

we can define a new « invariant » measure  $m_{\mu}$  on the  $\sigma$ -ring (<sup>6</sup>) L of Lebesgue measurable sets of R by

$$m_{\mu}(\Lambda) = \iint_{\Lambda} \mu(x, y) \, dx, \, dy$$

for  $\Lambda \in L$ . Then the local group of  $C^{(\alpha]}$ -homeomorphisms  $\{T_t\}$  is a  $m_{\mu}$ -measure true stream. Here  $m_{\mu}$  is a non-negative, completely additive set function on L. For any compact set  $K \subset R$ ,  $m_{\mu}(K) < \infty$ .  $m_{\mu}$  is an absolutely continuous measure in terms of the Lebesgue measure  $m_L$ , that is, if  $m_L(\Lambda) = 0$ , then  $m_{\mu}(\Lambda) = 0$ , The purpose of this note is to point out that if  $\mu(x, y) \ge 0$  and vanishes on no open set, even if  $\mu(x, y) \in C^{(\infty)}$ , there can exist a set  $Z \in L$  such that  $m_{\mu}(Z) = 0$  but  $m_L(Z) > 0$ . In this case the invariant measure  $m_{\mu}$ is not complete since there exist non-measurable subsets of Z. However in case  $\mu(x, y)$  is analytic in R,  $m_{\mu}(\Lambda) = 0$  if and only if  $m_L(\Lambda) = 0$ and  $m_{\mu}$  is a complete measure.

#### II. – Homeomorphisms which preserve measurability.

Any simply-connected plane region R can be mapped by a  $C^{(A)}$ -homeomorphism onto the entire plane. Thus we can replace R by the entire plane in any measure-theoretic considerations which are invariant under such a map.

Let  $\theta$  be a Borel set of Euclidean *n*-space. Let S be the  $\sigma$ -ring of

<sup>(\*)</sup> H. POINCARÉ, Méthodes nouvelles de la Mécanique céleste, t. III, chap. 26.

<sup>(6)</sup> Sometimes called  $\sigma$ -field. For notation, see P. HALMOS, Measure Theory (New-York, 1950).

all subsets of  $\theta$ . S actually forms an algebraic commutative Boolean (idempotent) ring with unity element under the operations of symetric difference  $\Delta$  for « addition » and intersection  $\Lambda$  for « multiplication ». Because of the simple, well-known formulae for union, difference, and complementation

(4) 
$$\begin{cases} \varphi \cup \psi = (\varphi \triangle \psi) \triangle (\varphi \cap \psi), \\ \varphi - \psi = \varphi \triangle (\varphi \cap \psi), \\ \varphi' = \varphi \triangle \theta \end{cases}$$

all of the set-theoretic properties of S are determined by its structure as an algebraic ring under  $\triangle$  and  $\bigcap$ . Let B be the  $\sigma$ -ring of Borel sets of  $\theta$  (smallest  $\sigma$ -ring containing the open sets of  $\theta$ ), with the Borel measure  $m_{\rm B}$ . B is an algebraic subring of S. The subset  $b \subset B$  of Borel sets with measure zero forms an ideal in the ring B. The  $\sigma$ -ring L of Lebesgue measurable subsets of  $\theta$  (all sets of the form  $\eta \bigtriangleup \nu$  where  $\eta \in B$  and  $\nu$  is a subset of some set in b) with Lebesgue measure  $m_{\rm (L)}$  is an algebraic subring of S and a superring of B. The ideal  $l \subset L$  of sets with Lebesgue measure zero is a superring of b, and consists of all sets of the form  $\beta \bigtriangleup \nu$  where  $\beta \in b$ and  $\nu$  is a subset of some set of b.

Let  $\theta_1$  and  $\theta_2$  be two Borel sets of Euclidean spaces and let  $S_1$ ,  $B_1$ ,  $b_1$ ,  $L_1$ ,  $l_1$  and  $S_2$ ,  $B_2$ ,  $b_2$ ,  $L_2$ ,  $l_2$  be the corresponding  $\sigma$ -rings described above. Let T be a homeomorphism of  $\theta_1$  onto  $\theta_2$ . Since T is a one-to-one point transformation, T induces a set mapping  $\tau$  from  $S_1$  onto  $S_2$ . Clearly  $\tau$  is an algebraic isomorphism of  $S_1$  onto  $S_2$ .

THEOREM I. — Let T be a homeomorphism of  $\theta_1$  onto  $\theta_2$ , Borel sets of Euclidean spaces. Then the induced set-map  $\tau$  is an isomorphism of  $B_1$  onto  $B_2$ . Furthermore the following four conditions are equivalent:

1. T is a  $m_{\rm L}$ -measurability preserving transformation;

2.  $\tau$  induces an isomorphism of L<sub>1</sub> onto L<sub>2</sub>;

3.  $\tau$  induces an isomorphism of  $b_1$  onto  $b_2$ ;

4.  $\tau$  induces an isomorphism of  $l_1$  onto  $l_2$ .

*Proof.* — Since T is a homoemorphism, open sets of  $\theta_1$  correspond to open sets of  $\theta_2$  under  $\tau$  and thus there is a one-to-one corres-

pondence between the sets of  $B_1$  and  $B_2$ . Since  $\tau$  is an isomorphism of  $S_1$  onto  $S_2$ ,  $\tau$  is therefore an isomorphism of  $B_1$  onto  $B_2$ .

Since  $\tau$  is an isomorphism of S<sub>1</sub> onto S<sub>2</sub>, conditions 1 and 2 are equivalent by the definition of measurability preserving transformations.

Suppose condition 3 holds; we shall deduce condition 2. If  $\Lambda_1 \in L_1$ , then  $\Lambda_1 = \eta_1 \bigtriangleup \nu_1$  where  $\eta_1 \in B_1$  and  $\nu_1 \subset \beta_1 \in b_1$ . Then

$$\tau(\Lambda_1) = \Lambda_2 = \tau(\eta_1) \bigtriangleup \tau(\nu_1) = \eta_2 \bigtriangleup \nu_2,$$

with ·

$$\tau(\eta_1) = \eta_2 \in \mathcal{B}_2$$
 and  $\tau(\nu_1) = \nu_2 \subset \tau(\beta_1) = \beta_2 \in b_2$ .

Thus  $\tau(\Lambda_1) = \Lambda_2 \in L_2$  and  $\tau$  maps  $L_1$  into  $L_2$ . Using the inverse homeomorphism  $T^{-1}$  we see that  $\tau^{-1}$  maps  $L_2$  into  $L_4$  and since  $\tau$  is one-to-one on  $S_1$ ,  $\tau$  is an isomorphism of  $L_1$  onto  $L_2$ .

Conversely suppose  $\tau$  is an isomorphism of  $L_1$  onto  $L_2$ . Suppose  $\tau(\beta_4) = \beta_2 \in B_2 - b_2$  for some  $\beta_4 \in b_4$ . Then  $m_B(\beta_2) > 0$ and there exists a set (<sup>7</sup>).  $\chi_{12} \subset \beta_2$  such that  $\chi_2 \in S_2 - L_2$ . But then  $\tau^{-4}(\chi_2) = \chi_4 \subset \beta_4$  and thus  $\chi_4 \in L_4$ . But  $\tau(\chi_4) \notin L_2$  and this contradicts the hypothesis that  $\tau$  maps  $L_4$  onto  $L_2$ . Thus  $\tau$  maps  $b_4$ into  $b_2$  and, as above,  $\tau$  is an isomorphism of  $b_4$  onto  $b_2$ .

Since  $b_i = B_i \cap l_i$  (i=1,2), condition 4 clearly implies 3. Because  $l_i$  consists of those sets of  $L_i$  which are contained in sets of  $b_i$  (i=1,2), conditions 2 and 3 imply 4.

COROLLARY. — If T is a C<sup>(1)</sup>-homeomorphism of  $\theta_1$  onto  $\theta_2$ , open sets of an Euclidean n-space, then T is  $m_{L}$ -measurability preserving.

*Proof.* — We need show only that  $\tau(b_1) \subset b_2$ . If  $\beta_i \in b_i$ , then  $\beta_i$  has a covering by a countable number of compact closed *n*-balls  $\overline{R}_i$  with

$$\sum_{i=1}^{\infty} m_{\rm B}(\bar{\rm R}_i) \equiv 1.$$

In  $\overline{R}_i$  the continuous Jacobian of T is bounded,  $|J(T)| < K_i$  Then

$$n_{\mathbf{B}}(\beta_1 \cap \overline{\mathbf{R}}_i) \equiv 0$$

<sup>(7)</sup> See P. HALMOS, op. cit., p. 70. The existence of a non-measurable linear set combined with an application of Fubini's theorem gives this result.

and we cover  $\beta_i \cap \overline{R}_i$  by an open set  $o_i$  with  $m_B(o_i) = \frac{\delta}{2K_i}$ . Then

$$m_{\mathbf{B}}(\tau(\mathbf{o}_i)) \leq \frac{\delta}{2^i}$$

and thus

$$\sum_{i=1}^{\infty} m_{\mathrm{B}}(\tau(\mathbf{o}_{t})) \leq \delta.$$

Thus  $m_{B}(\tau(\beta_{1})) \leq \delta$  and therefore, since  $\delta$  is an arbitrary positive number,  $\tau(\beta_{1}) = \beta_{2} \in l_{2}$ . But certainly  $\beta_{2} \in B_{2}$ . Thus  $\beta_{2} \in b_{2}$ . Therefore  $\tau(b_{1}) \subset b_{2}$  and using  $T^{-1}$ , which also has a continuous Jacobian,  $\tau^{-1}(b_{2}) \subset b_{1}$  and therefore  $\tau$  maps  $b_{1}$  onto  $b_{2}$ .

An example of a  $C^{(0)}$ -homeomorphism of the linear interval [0, 1] onto itself which is not measurability preserving is given by T

Q. E. D.

$$x \rightarrow f(x) = \frac{1}{2} [x + \psi(x)]$$

where  $\psi(x)$  is the continuous, but not absolutely continuous, nondecreasing, Cantor function (\*). If K is the compact Cantor set, then  $m_{\rm B}({\rm K}) = 0$  but  $m_{\rm B}(f({\rm K})) = \frac{1}{2}$ .

#### III. - Completeness of invariant measures.

We shall construct a function  $\mu(x, y)$  such that :

- 1.  $\mu(x, y) \in C^{(\infty)}$  in the entire plane;
- 2.  $\mu(x, y) \ge 0$  and the plane set Z defined by  $\mu(x, y) = 0$  contains no interior points;
- 3.  $m_{\rm B}({\rm Z}) = \infty$ .

Let z be a compact, nowhere dense point set of the real line such that the linear Borel measure  $m'_{\rm B}(z) > 0$ . A closed point set is nowhere dense if and only if it contains no interior point. Such a

<sup>(8)</sup> P. HALMOS, op. cit., p. 83.

set z is given in Hobson ( $^{9}$ ) and can be described briefly as follows.

Let [0, 1] be divided into m > 2 equal parts and the last exempted from further division. Then let the remaining m - 1 parts each be divided into  $m^2$  equal parts, the last of each being exempted from further division. Let the remaining parts be then divided into  $m^3$ equal parts, the last of these in each case being exempted from further division. If this process is carried out a countable number of times, the endpoints of the divisions, together with their accumulation points, form a nowhere dense, compact set z.

This set z has mesure

$$0 < \prod_{i=1}^{\infty} \left( \mathbf{I} - \frac{\mathbf{I}}{m^{i}} \right) <$$

I

for after i operations, the measure of the union of the exempted segments is

$$\frac{1}{m} + \frac{m-1}{m^3} + \frac{(m-1)(m^2-1)}{m^6} + \ldots + \frac{(m-1)(m^2-1)\ldots(m^{i-1}-1)}{m^{\frac{i(i+1)}{2}}}$$

or

$$\mathbf{I} - \left(\mathbf{I} - \frac{\mathbf{I}}{m}\right) \left(\mathbf{I} - \frac{\mathbf{I}}{m^2}\right) \cdots \left(\mathbf{I} - \frac{\mathbf{I}}{m^i}\right).$$

Thus  $m'_{B}(z) > 0$  and indeed, may be chosen arbitrarily close to one if the dividing ratio *m* be large enough. The complement of z in [0, 1] is a countable union of disjoint open intervals  $I_n$  and zconsists of the closure of the endpoints of the  $I_n$ .

The distance from a real number x to the compact set z is  $d(x) \ge 0$ and d(x) = 0 if and only if  $x \in z$ .

Then for any point (x, y) of the plane we define

	$(\mu(x, y) = 0$	in case $x \in z$ .
•	$\mu(x, y) = e^{-\frac{1}{x^2}}$	in case $x < 0$ .
(5)	$\mu(x, y) = e^{-\frac{1}{(x-1)^2}}$	in case $x > 1$ .
	$\mu(x, y) = e^{-\frac{\overline{\Delta_n^2}}{4} - \left(\frac{\overline{\Delta_n}}{2} - d(x)\right)^2}$	in case $x \in I_n$

(\*) A. E. HOBSON, The theory of functions a real variable (Cambridge, 1921), p. 119 and also p. 164.

where  $A_n = b_n - a_n$  is the length of  $I_n = (a_n, b_n)$ . Then  $\mu(x, y) \ge 0$ and, defining the plane set  $(x, y) \in \mathbb{Z}$  in case  $x \in z$ ,  $\mu(x, y) = 0$  if and only if  $(x, y) \in \mathbb{Z}$ . The set Z is closed and contains no interior points of the plane. Moreover  $m_B'(\mathbb{Z}) = \infty$ . We shall show that  $\mu(x, y) \in \mathbb{C}^{(\infty)}$  and then clearly

$$m_{\mu}(\mathbf{Z}) = \iint_{\mathbf{Z}} \mu(x, y) \, dx \, dy = 0.$$

If x < 0 or x > 1, then  $\mu(x, y) \in C^{(\infty)}$ . If  $x \in I_n$ , that is,  $a_n < x < b_n$ , then

$$\frac{\Delta_n^2}{4} - \left(\frac{\Delta_n}{2} - d(x)\right)^2 = (x - a_n) (b_n - x)$$

and thus  $\mu(x, y) \in C^{(\infty)}$  for these regions.

For  $x \in I_n$ ,

$$\frac{\Delta_n^2}{4} - \left(\frac{\Delta_n}{2} - d(x)\right)^2 = 2 \,\Delta_n d - d^2 < d - d^2 < d.$$

Thus

$$|\mu(x, \gamma)| \leq e^{-\frac{1}{d(x)}} < \frac{1}{d(x)^2} e^{-\frac{1}{d(x)}}.$$

By induction it is easy to show that

$$\frac{\partial^{k} \mu}{\partial x^{k}} = \frac{\mathbf{P}_{k}((x-a_{n})), (b_{n}-x)}{[(x-a_{n})(b_{n}-x)]^{2^{k}}} e^{-\frac{1}{(x)-a_{n}(b_{n}-x)}}$$

where  $P_k$  is a polynomial and since

$$d(x) \leq |x-a_n| < \mathfrak{l}, \quad d(x) \leq |b_n-x| < \mathfrak{l},$$

we have

$$|\mathbf{P}_k((x-a_n)(b_n-x))| < \overline{\mathbf{P}}_k$$

an upper bound independent of  $I_n$  and thus

(6) 
$$\left|\frac{\partial^k \mu}{\partial x^k}\right| < \frac{\overline{\mathrm{P}}_k}{d(x)^{2^{k+1}}} e^{-\frac{1}{d(x)}}$$
 for  $k = 0, 1, 2, \ldots$ 

For sufficiently small positive d, the function  $\frac{1}{d^{2^{k+1}}}e^{-\frac{1}{d}}$  is strictly increasing.

Finally consider a point  $(x_0, y_0) \in \mathbb{Z}$ . Then  $\mu(x_0, y_0) = 0$  and

for  $|x-x_0| < \delta$ , a sufficiently small positive number,  $d(x) < \delta$  and

$$p \leq |\mu(x, y) - \mu(x_0, y_0)| \leq \frac{\overline{\overline{P}}_0}{d(x)^2} e^{-\frac{1}{d(x)}} < \frac{\overline{\overline{P}}_0}{\delta^2} e^{-\frac{1}{\delta}}.$$

Thus  $\mu(x, y)$  is continuous at each point of Z and  $\mu(x, y) \in C^{(0)}$  in the plane. We prove  $\mu(x, y) \in C^{(\infty)}$  by induction. Suppose

$$\mu(x, y) \in \mathbb{C}^{(k)}$$
 and  $\frac{\partial^k \mu}{\partial x^k}(x_0, y_0) = 0.$ 

Then for  $0 < |x - x_0| < \delta$ , a sufficiently small positive number,  $d(x) < \delta$  and

$$\frac{\frac{\partial^k \mu(x, y)}{\partial x^k} - \frac{\partial^k \mu(x_0, y_0)}{\partial x^k}}{x - x_0} \bigg| < \frac{\overline{\mathrm{P}_k e^{-\frac{1}{\delta}}}}{\delta \delta^{z^{k+1}}}.$$

Thus  $\frac{\partial^{k+1}\mu(\gamma_0, x_0)}{\partial x^{k+1}} = 0$ . From the inequality  $\left| \frac{\partial^{k+1}\mu}{\partial x^{k+1}} \right| < \frac{\overline{P}_{k+1}}{d(x)^{2^{k+2}}} e^{-\frac{1}{d(x)}}$ 

it is clear that

$$\lim_{\substack{x=x_0\\y=y_0}}\frac{\partial^{k+1}u(x,y)}{\partial x^{k+1}}=0$$

and thus  $\frac{\partial^{k+1}\mu(x, y)}{\partial x^{k+1}}$  exists and is continuous at each point of Z and therefore at every point of the plane. Since  $\frac{\partial \mu}{\partial y} \equiv 0$  we have

$$\mu(x, y) \in \mathbb{C}^{(k+1)}$$
 and  $\frac{\partial^{k+1} \mu(x_0, y_0)}{\partial x^{k+1}} = 0.$ 

Therefore the induction is complete and  $\mu(x, y) \in C^{(\infty)}$  in the plane.

THEOREM II. — Let R be simply-connected plane region. Then there exists a real, first oder non-singular ordinary differential equation.

(1) 
$$\frac{dx}{dt} = f(x, y), \quad \frac{dy}{dt} = g(x, y) \quad (f^2 + g^2 > 0)$$

with  $f, g \in C^{(\infty)}$  in R, with an integrating factor  $\mu(x, y) \in C^{(\infty)}$ , defining an invariant measure

$$m_{\mu}(\Lambda) = \iint_{\Lambda} \mu(x, y) \, dx \, dy$$

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on the plane Lebesgue  $m_{\rm L}$ -measurable sets, such that

1.  $\mu(x, y) \in C^{(\infty)}$  in R; 2.  $\mu(x, y) \ge 0$ ; 3.  $\mu(x, y)$  vanishes on no open set; 4.  $\frac{\partial}{\partial x}(\mu f) + \frac{\partial}{\partial y}(\mu, g) \equiv 0$  in R

and the  $m_{\mu}$ -measure is not complete.

If f, g are only in  $C^{(1)}$  with  $f^2 + g^2 > 0$  in R, and if  $\mu(x, y) \in C^{(A)}$ in R is an integrating factor, that is, satisfies conditions 2, 3, and 4, then the  $m_{\mu}$ -measure is complete, that is,  $m_{\mu}(\Lambda) = 0$  if and only if  $m_{L}(\Lambda) = 0$ .

Proof. — We first prove the theorem in the case R is the entire plane and then use theorem I for the general simply-connected region. Consider the differential system

(7) 
$$\frac{dx}{dt} = \mu(x, y), \qquad \frac{dy}{dt} = -2 \frac{\partial \mu}{\partial x} y + 1$$

where  $\mu(x, y)$  is the function described above. This system is  $C^{(\infty)}$ and also non-singular since  $\mu(x_0, y_0) = 0$  if and only if  $(x_0, y_0) \in \mathbb{Z}$ in which case  $\frac{\delta\mu}{\delta x} = 0$  and  $(-2\mu_x y + 1) = 1$ . Moreover, it is not exact but has an integrating factor of  $\mu(x, y)$  since

$$\frac{\partial}{\partial x}(\mu^2) + \frac{\partial}{\partial y}(-2\mu_x\mu_y + \mu) \equiv 0.$$

To show the  $m_{\mu}$ -measure is not complete let  $\chi \in \mathbb{Z}$  be a non $m_{L}$ -measurable set. This is possible since  $m_{L}(\mathbb{Z}) = \infty$ . But

$$m_{\mu}(\mathbf{Z}) = \iint_{\mathbf{Z}} \mu(x, y) \, dx \, dy = 0$$

and thus  $\chi$  is a non- $m_{\mu}$ -measurable subset of a set Z of  $m_{\mu}$ -measure zero. Therefore the  $m_{\mu}$ -measure is not complete.

Next let  $\mu(x, y) \in C^{(A)}$  be any integrating factor of 1 in the plane. Clearly if

 $m_{\rm L}(\Lambda) \equiv 0$ 

then

$$m_{\mu}(\Lambda) = \iint_{\Lambda} \mu \, dx \, dy = 0.$$

Conversely let  $m_{\mu}(\Lambda) = 0$ . Then there exist two sets  $\Lambda_1$ ,  $\Lambda_2 \in L$  such that  $m_{L}(\Lambda_1) = 0$ , and

$$\Lambda_1 \cap \Lambda_2 = 0, \quad \Lambda_1 \cup \Lambda_2 = \Lambda, \quad \text{and} \quad \Lambda_2 \subset \mathbb{Z},$$

the set of zeros of  $\mu(x, y)$ . We shall show that  $m_{\rm L}({\rm Z}) = 0$ .

Since Z is closed,  $Z \in L$ . Then for almost all horizontal lines  $h_c: y = y_c$ , the linear measure of  $Z \cap h_c$  exists. If  $Z \cap h_c$  contains only a countable number of points, it has linear measure zero. If  $Z \cap h_c$  contains a non-countable number of points, then since  $\mu(x, y_c)$  is an analytic function of one argument,  $Z \cap h_c = h_c$ and thus has infinite linear measure. However there are only a countable number of lines  $h_c$  on which  $m'_L(Z \cap h_c) = \infty$ . For if otherwise, there would be a finite accumulation point  $\overline{y}$  of the corresponding ordinates and then, for each fixed  $x_0$ ,  $\mu(x_0, y)$  is an analytic function of y and must vanish. Thus  $\mu(x, y) \equiv 0$  which contradicts the hypothesis of the theorem that  $\mu(x, y)$  vanishes on no open set. Therefore on almost all horizontal lines  $m'_L(Z \cap h_c) = 0$ . Therefore by Fubini's theorem  $m_L(Z) = 0$ . Since  $\Lambda \subset \Lambda_1 \cup Z$ and  $m_L(\Lambda_1) = 0$ , we have  $m_L(\Lambda) = 0$ .

Next consider any subset  $\Lambda' \subset \Lambda$ . Then  $\Lambda' \in l$  and  $m_{\rm L}(\Lambda') = 0$ . Thus  $m_{\mu}(\Lambda') = 0$  and the  $m_{\mu}$ -measure is complete. The theorem is proved in case R is the entire plane.

Again return to the case of a general simply-connected plane region R. Let  $T: (x, y) \rightarrow (u(x, y), v(x, y))$  be a  $C^{(A)}$ -homeomorphism (<sup>10</sup>) of the plane onto the (u, v)-region R. By T the

(10) First map the plane onto the square  $|x^1| < 1$ ,  $|y^1| < 1$  by

$$x^1 = \frac{2}{\pi} \operatorname{tg}^{-1} x, \qquad y^1 = \frac{2}{\pi} \operatorname{tg}^{-1} y.$$

Then use the Riemann conformal mapping theorem.

differential system (7) is carried into a non-singular C<sup>( $\infty$ )</sup>-differential system defined in R

(8) 
$$\frac{du}{dt} = u_x f + u_y g = \mathbf{F}(u, v), \qquad \frac{dv}{dt} = v_x f + v_y g = \mathbf{G}(u, v)$$

where F,  $G \in C^{(\infty)}$  in R and  $F^2 + G^2 > 0$ . The corresponding integrating factor is

$$\mathbf{M}(u, v) = \mu \left| \frac{\partial(x, y)}{\partial(u, v)} \right|.$$

with the corresponding integral  $\psi(x(u, v), y(u, v))$ , as is seen from the following matrix equation

(9) 
$$\begin{pmatrix} \psi_{\mu} \\ \psi_{\nu} \end{pmatrix} = \begin{pmatrix} x_{\mu} & y_{\mu} \\ x_{\nu} & y_{\nu} \end{pmatrix} \begin{pmatrix} \psi_{x} \\ \psi_{y} \end{pmatrix} = \begin{pmatrix} x_{\mu} & y_{\mu} \\ x_{\nu} & y_{\nu} \end{pmatrix} \mu \begin{pmatrix} -g \\ f \end{pmatrix}$$

and

$$\begin{pmatrix} -g \\ f \end{pmatrix} = \begin{pmatrix} y_{\nu} & -y_{u} \\ x_{\nu} & x_{u} \end{pmatrix} \begin{pmatrix} -G \\ F \end{pmatrix}.$$

Thus

(10) 
$$\begin{pmatrix} \psi_u \\ \psi_v \end{pmatrix} = \mu(x(u, v), y(u, v)) \frac{\partial(x, y)}{\partial(u, v)} \begin{pmatrix} -G \\ F \end{pmatrix} = \pm M \begin{pmatrix} -G \\ F \end{pmatrix}.$$

Clearly  $M(u, v) \geq o$ ,  $M \in C^{(\infty)}$  and vanishes only on the set T(Z) which contains no interior points. Since T is a  $C^{(4)}$ -homeomorphism it preserves the set properties measurability and also of having zero measure. Thus  $m_L(T(Z)) > o$  and thus the  $m_M$ -measure is not complete.

A similar argument proves that, in case  $\mu(x, \gamma) \in C^{(A)}$  is an integrating factor of (1) in R, a simply-connected region, the resulting invariant measure is complete. Q. E. D.

COROLLARY. - Let

(1) 
$$\frac{dx}{dt} = f(x, y) \qquad \frac{dy}{dt} = g(x, y) \qquad (f^2 + g^2 > 0)$$

with  $f, g \in C^{(A)}$  in a simply-connected plane region R have a principal integral  $\psi(x, y)$  such that

1. 
$$\psi(x, y) \in C^{(A)}$$
 in R;

.

2.  $\psi(x, y)$  is constant on no open set of R;

3.  $\psi(x, y)$  is constant along each solution curve of (1).

Then the differential system (1) has a complete invariant measure.

Proof. — Let

$$\mu(x, y) = \left[\frac{\psi_x^2 + \psi_y^2}{f^2 + g^2}\right]^{\frac{1}{2}}$$

be the integrating factor corresponding to the integral  $\psi(x, y) \in C^{(A)}$ . Then  $\mu(x, y) \in C^{(A)}$  and, by the theorem, the invariant measure

$$m_{\mu}(\Lambda) = \iint_{\Lambda} \mu(x, y) \, dx \, dy$$

is complete.

Q. E. D.

We shall discuss criteria that (1) should have an analytic integral in a later paper.

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