REVISTA DE LA UNION MATEMATICA ARGENTINA

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VOLUMEN 24, NUMEROS 2 Y 3 1969

BAHIA BLANCA 1969

UNION MATEMATICA ARGENTINA

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REVISTA SEMESTRAL

REVISTA DE LA UNION MATEMATICA ARGENTINA Volumen 24, Número 2, 1968

INDUCED SHEAVES AND GROTHENDIECK TOPOLOGIES by Juan José Martínez

INTRODUCTION. The theory of sheaves, as it is exposed in the classical book of R. Godement {2}, has been generalized in successive stops. Ending this process, M. Artin introduced the notion of Grothendieck topology and developed the fundamental part of the theory in a functorial way (cf. {1}). Although, the concept of Grothen dieck topology seems to be insufficient to relate certain aspects of the theory of sheaves; for example, the notion of subspace (not necessarily open !) is omited and so, induced sheaves and relative cohomology must be ignored.

The purpose of this paper is to obtain the essential results about induced sheaves (the concept of topological category enable us to work in this direction; cf. \$1). Topological methods play an important role in the problems in question, as Godement shows (cf. $\{2\}$, Ch. II, \$2.9). Therefore, we are forced to introduce a various kind of axioms, valids -of course- in the classical situation of a topological space. We mention that the results of this paper are useful also in not conventional cases, namely, the "étale" Grothendieck topology for preschemes (cf. $\{1\}$, Ch. III).

Results and notations of Artin's seminar ({1}, Ch. I, II) are continuosly used, frequently without specific reference. This results are stated in {1}for sheaves of abelian groups, but all of them could be generalized taking an arbitrary category of values and inserting axioms where necessary.Here, we have followed the abstract formulation (the basic facts about limit of functors, existence of injectives, adjoint situations, derived functors of a composition, etc. are stated in the usual literature; for example, cf. $\{3\}$). Of course, the reader could suppose that all sheaves in this paper are abelian sheaves.

1. TOPOLOGICAL CATEGORIES AND INDUCED SHEAVES. This section is of introductory character. Its aim is to lay down the terminology used throughout this paper and to collect the basic facts. We begin with the following:

DEFINITION 1.1. A topological category is a triple (M,T,ϕ) such that M is a category, T is a family of Grothendieck topologies $(T_M)_{M \in ObM}$ wis a family of morphisms in M $(\phi_M: M \to X)_{M \in ObM}$ and the following axioms are satisfied:

tcl) For all object M in M Cat T_M is a full subcategory of M. X is an object of Cat T_X and ϕ_X is the identity morphism e_X of X. tc2) The diagram



is commutative, for all morphism $f: U \rightarrow V$ in Cat T_{y} .

tc3) M has fibered products of the form $\underset{\Psi_{\rm U}}{U \times M}$ (briefly noted U $\times_{\rm X}$ M), where U is an object of Cat T $_{\rm X}$ and M is an object of M, such that

 $U \in Ob(Cat T_x) \Longrightarrow U \times_X M \in Ob(Cat T_M)$

 $(f_i: U_i \to U)_{i \in I} \in Cov \ T_X \xrightarrow{\longrightarrow} (f_i \times_X e_M: U_i \times_X M \to U \times_X M)_{i \in I} \in Cov \ T_M$ for all object M of M.

REMARKS 1.2. i) Axioms tc1 and tc2 tell us that M is an object of Cat T_M , because X $\times_X M$ = M and X is an object of Cat T_X .

ii) Recall that if X is an object in a category M, then is called prefinal (resp. final) iff $\operatorname{Hom}_{M}(M,X) \neq \emptyset$ (resp. $\operatorname{Hom}_{M}(M,X)$ is a set of one element), for all M ε ObM. If < X > is the discret subcatego ry of M associated to X, one easily checks that the following state ments are equivalent (cf. {1} Ch. I, §0):

a) M satisfies axiom L1 and < X > is a final subcategory of M.

b) X is a final object of M.

c) M satisfies axiom L2 and X is a prefinal object of M such that $Hom_M(X,X) = \{e_x\}$.

Clearly, if (M,T,ϕ) is a tc (topological category) then X is a prefinal object of M.

Let M be a category with final object X and let ϕ be the family of morphisms canonically associated to X. If T is a family of topologies satisfying tc1 and tc3, respect to M and ϕ , then (M,T, ϕ) is a tc of the following type:

DEFINITION 1.3. A topological category (M,T,ϕ) is called tc⁰ iff it satisfies:

tc2') For all morphism $f: M \rightarrow N$ in M the diagram

 $\begin{array}{c} M \xrightarrow{\mathbf{f}} N \\ \phi_{\mathbf{M}} & & \phi_{\mathbf{N}} \end{array}$

is commutative.

Given a tc C = (M,T, ϕ), we shall be using a naive nomenclature: M is called the category of subspaces of C (consistently, an object M in M is called a subspace of C). The object X is referred to as the space of C and so, an object M of M is also called a subspace of X refering to ϕ_M as the inclusion morphism of M in X.

If M is a subspace of M, T_{M} is called the relative topology of M and Cat T_{M} is called the category of relative open objects of M. Abusing language, T_{X} is called the topology of C and Cat T_{X} is called the category of open subspaces of C.

A morphism f: M \longrightarrow N in M is called a ϕ -morphism iff the diagram



is commutative. We define a category ${\rm M}_{_{\rm A}}$ putting:

ObM = ObM

Hom M_{ϕ} : ϕ -morphisms of M

 M_{ϕ} is a subcategory of M and clearly is a full subcategory (equivalently, is equal to M) iff C is tc⁰.

DEFINITION 1.4. A morphism of topological categories F: C \longrightarrow C' is a functor F: M \longrightarrow M' such that:

mtcl) For all object M of M F/Cat T_{M} is a morphism of topologies, of T_{M} in $T'_{F(M)}$.

mtc2) $F(\phi) = \phi'$ (i.e. F(X) = X' and $F(\phi_M) = \phi'_{F(MY)}$, for all MeObM). mtc3) F preserves fibered products of the form $U \times_X M$, UeOb(Cat T_X) and $M \in ObM$.

REMARKS 1.5. i) Now we can talk about the category of small topological categories.

ii) If F: C \longrightarrow C' is a mtc then we have

 $U \in Ob(Cat T_{X}) \longrightarrow F(U) \times_{X} F(M) \in Ob(Cat T'_{F(M)})$ $(f_{i}: U_{i} \longrightarrow U)_{i \in I} \in Cov T_{X} \longrightarrow (F(f_{i}) \times_{X} e_{F(M)}: F(U_{i}) \times_{X} F(M) \longrightarrow$ $+ F(U) \times_{X} F(M)_{i \in I} \in Cov T'_{F(M)}$ for all subspace M of C.

iii) A mtc F: C \longrightarrow C' induces for each subspace M of C a mor phism of topologies F/Cat $T_M: T_M \longrightarrow T'_{F(M)}$ and so, induces the usual functors (direct and inverse image) between the corresponding categories of presheaves or of sheaves.

Let M be a subspace of a tc C. If A is an arbitrary category, the category of presheaves $P(T_M, A)$ is briefly denoted by P_M , and a presheaf in P_M is called a presheaf over M. Similarly, if A is a category with products, S_M denotes the category of sheaves $S(T_M, A)$, and a sheaf in S_M is called a sheaf over M.

 $\rho_M: T_X \longrightarrow T_M$ is the morphism of topologies defined by the assignment of objects U \longrightarrow U $\times_X M$.

A category A will be called:

o) A0 iff it is a complete category (respect to functorial direct limits) with products and zero object.

i) A1 iff it is A0 and abelian.

ii) A2 iff it is A1 and satisfies the Grothendieck axiom A.B.5.iii) A3 iff it is A2 and has a generator.

Let C be a tc and let A be an A1 category (as category of values).

DEFINITION 1.6. If M is a subspace of C and F is a sheaf over X, then we call ρ_{M_S} (F) the sheaf induced by F over M, and we denote it by F/M.

DEFINITION 1.7. If M is a subspace of C and $\alpha: F \longrightarrow G$ is a morphism of sheaves over X, then we call $\rho_{M_S}(\alpha)$ the morphism induced by α over M, and we denote it by α/M .

REMARKS 1.8. i) Since $\rho_X = e_{T_X}$ it is clear that F/X = F and $\alpha/X = \alpha$ ii) Since ρ_{M_S} is a functor it is clear also that $e_F/M = e_{F/M}$ and $(\beta\alpha)/M = (\beta/M)(\alpha/M)$.

iii) Remark that expressions of the type (F/M)/N have no sense here, because the "absolute" topology $T_{\rm X}$ plays a special role in our devel opments.

Now we need to prove some previous results. In the next lemma, and only in the next, A may be an arbitrary category.

LEMMA 1.9. If f: $K \longrightarrow K'$ is a morphism of small categories, the following statements are true:

i) If f is a full and representative functor, then $f^P: P' \longrightarrow P$ is full.

ii) If f is representative and $\alpha \in$ Hom P' is such that $f^P(\alpha)$ is an isomorphism, then α is also an isomorphism.

Proof i) We want to show that the function $\operatorname{Hom}_{p}(P_{1}, P_{2}) \longrightarrow \longrightarrow \operatorname{Hom}_{p}(f^{p}(P_{1}), f^{p}(P_{2}))$, $\alpha \longrightarrow f^{p}(\alpha)$, is surjective. Given $\beta \in \operatorname{Hom}_{p}(f^{p}(P_{1}), f^{p}(P_{2}))$ we define a morphism $\alpha \in \operatorname{Hom}_{p}(P_{1}, P_{2})$ in the following way: since f is representative, given an object V in K' there exists an object U in K such that f(U) = V; therefore, we take $\alpha(V) = \beta(U)$. The good definition of α is obtained by the following argument: if U' is an object of K such that f(U')=V, since f is a full functor there exists a morphism m: U' \longrightarrow U such that $f(m) = e_{V}$. Now, since β is a morphism of presheaves, we have the commutative diagram

$$f^{p}(P_{1})(U) \xrightarrow{\beta(U)} f^{p}(P_{2})(U)$$

$$f^{p}(P_{1})(m) \xrightarrow{f^{p}(P_{1})(U')} f^{p}(P_{2})(U')$$

i.e. we have the commutative diagram

and so, $\beta(U) = \beta(U')$. It is trivial that $f^{P}(\alpha) = \beta$. ii) Given $\alpha \in \text{Hom P'}$ and $V \in \text{ObK'}$, observe that $V = f(U) \longrightarrow \alpha(V) = f^{P}(\alpha)(U)$.

COROLLARY 1.10. If $f: T \longrightarrow T'$ is a morphism of topologies, the following statements are true:

i) If f is a full and representative functor, then $f^{S}\colon S' \longrightarrow S$ is full.

ii) If f is representative and a ϵ Hom S' is such that $f^{\bf s}(\alpha)$ is an isomorphism, then α is also an isomorphism.

Proof: Apply the lemma, taking in mind that $if^s \approx f^p i'$, where i (resp. i') is the inclusion functor of S (resp. S') in P (resp. P')

COROLLARY 1.11. If $f: T \longrightarrow T'$ is a full and representative mor - phism of topologies, then $f_s: S \longrightarrow S'$ is a representative functor.

Proof: Since f_s is left adjoint to f^s , there exists a canonical morphism of functors $\Lambda: f_s \circ f^s \longrightarrow e_{S'}$. Now, since f^s is full by 1.10,i, we have that $f^s(\Lambda_{F'}): f^s(f_s \circ f^s(F')) \longrightarrow f^s(F')$ is an isomorphism, for all sheaf F' in S'. Therefore, applying 1.10, ii, $\Lambda_{F'}: f_s f^s(F') \longrightarrow F'$ is also an isomorphism; and so, given a sheaf F' in S' the sheaf $f^s(F')$ is a preimage by f_s of F'.

The situation above suggest us the following

DEFINITION 1.12. A tc C is called tc^1 iff the morphism ρ_M is a full and representative functor, for all subspace M of C.

REMARK 1.13. If C is tc¹ we can apply both corollaries to the morphism $\rho_M: T_X \longrightarrow T_M$. In particular, 1.11 tell us that the restriction functor ./M: $S_X \longrightarrow S_M$ is representative, for all subspace M of a tc¹ C.

In order to obtain the classical theorems about " characteristic " sheaves (cf. {2}, Ch. II § 2.9) our first result is

LEMMA 1.14. Let M be a subspace of a tc¹ C and let F be a sheaf over X. If we define the sheaf F^{M} by $F^{M} = f^{s}f_{s}(F)$ then $F^{M}/M \simeq F/M$. (f is ρ_{M}).

Proof: Adjointness gives us a canonical morphism of functors $\Delta : e_{S}^{+} \rightarrow f^{s}f_{s}$; since f is a full functor (cf. 1.10,i) $f_{s}(\Delta_{F}): f_{s}(F) \rightarrow f_{s}(F^{M})$ is an isomorphism. Hence, $F/M \simeq F^{M}/M$.

The technique concerning to open subspaces will be obtained using the following type of categories

DEFINITION 1.15. A tc C is called tc^2 iff for all open subspace A of C the following conditions are satisfied:

i) T_A is a subtopology of T_X .

ii) If V is open in A, then $V \times_X A = V$ (i.e. $e_V \times_X \phi_A : V \times_X A \longrightarrow V \times_X X = V$ is an isomorphism).

LEMMA 1.16. If A is an open subspace of a tc² C and V is open in A, then < (V, h_V) > is initial in I_V^f , where f: $T_X \longrightarrow T_A$ is the morphism ρ_A and h_V : V \longrightarrow f(V) is the inverse morphism of $e_V x_X \phi_A$

Proof: Let (U,n) be any object in I_V^f . If p_U : f(U) \longrightarrow U denotes the first projection, we define a morphism m: V \longrightarrow U by m = p_U^n . We claim that m: (V,h_V) \longrightarrow (U,n) is a morphism in I_V^f ; to prove this we are reduced to check that the diagram



is commutative. In fact, if $g_V : f(V) \longrightarrow V$ is the morphism $e_V x_X \phi_A$ we have the commutative diagram

$$\begin{array}{ccc} f(V) & \xrightarrow{ng_V} & f(U) \\ g_V & \downarrow & & \downarrow & p_U \\ V & \xrightarrow{m} & U \end{array}$$

Therefore, recalling that g_V is the first projection of f(V), an uniqueness result on fibered products yields $f(m) = ng_V$. Hence, $f(m)h_V = n$.

The last thing to check is that $End((V,h_V))$ is a set of one ele ment. In fact, if r: $(V,h_V) \longrightarrow (V,h_V)$ is a morphism in I_V^f , then the diagram



73

is commutative, i.e. $f(r) = e_{f(v)}$. Therefore, since the diagram



is also commutative, results $r = e_v$.

COROLLARY 1.17. If A is a complete category, Cat $T_{\rm X}$ has fibered products and V is open in A, then the following statements are true:

i) If P is a presheaf over X, then $f_{p}(P)(V) \simeq P(V)$.

ii) If in addition A has products and F is a sheaf over X, then $f_pi(F)$ is a sheaf over $A(i: S_X \longrightarrow P_X$ is the inclusion functor). iii) If A is abelian too and F is a sheaf over X, then $F/A(V) \simeq F(V)$.

Proof: i) Since Cat T_X has fibered products and f preserves fibered products (because is a morphism of type ρ_M), the category I_V^f satisfies axiom L1* (cf. {1} II, Th. 4.14). Therefore, applying the lemma, we see that (V, h_V) is an initial object in I_V^f , and so $f_p(P)(V) = \underline{\lim} P_V \approx P_V((V, h_V)) = P(V)$.

ii) Applying i, check the definition of sheaf.

iii) It is clear by ii that $F/A \simeq f_p(F)$. Hence, i yields the desired result.

2. COMPLEMENTED TOPOLOGICAL CATEGORIES AND CLOSED SUBSPACES.

At this point, we need the notion of closed object in a topological category. Since we have the concept of open object, thinking in the closed sets of a topological space it is enough to find a notion replacing the set-theoretic operation of complement. Thus, we give the following

DEFINITION 2.1. A complemented topological category is a tc C, to gether with a functor c: $M^* \longrightarrow M$ such that, if $\theta = cX$ and $U_X = CatT_X$, the following axioms are satisfied:

ctcl) c is an involution functor (i.e. $c^* \circ c = e_{M^*}$).

ctc2) $\theta \in ObU_{X}$ and there exists $(U_{i} \longrightarrow \theta)_{i \in I} \in Cov T_{X}$ such that $I = \phi$

ctc3) $CM \times_{\mathbf{y}} M = \theta$ and $\theta \times_{\mathbf{y}} M = \theta$, for all $M \in ObM$.

REMARKS 2.2. i) $F_{\rm X} = c U_{\rm X}$ is called the category of closed subspaces of C, and c is called the complement operator of C.

ii) Axiom ctc3 says that the diagrams

$$\begin{array}{c} \theta \xrightarrow{c(\phi_{cM})} M \\ c(\phi_{M}) \downarrow & \downarrow \phi_{M} \\ cM \xrightarrow{\phi_{cM}} X \end{array} \qquad \qquad c(\phi_{X}) = e_{\theta} \downarrow & \downarrow \phi_{M} \\ \theta \xrightarrow{\phi_{\theta}} X \end{array}$$

are fibered products. We recall that (in the following proofs) we only need the first condition of ctc3 for closed subspaces, and the second for open subspaces.

iii) ctc1 tell us:

a) the complement of a closed object is open.

(The dual proposition is trivially true). From the definition of $\boldsymbol{\theta}$ and ctc2 we obtain:

b) θ is open and closed.

Therefore:

c) X is open and closed.

Using ctc2 and the second condition of ctc3, we see:

d) If M is a subspace of C, then $\theta \in Obl_{M}$ and there exists a covering $(V_{i} \longrightarrow \theta)_{i \in I}$ in T_{M} such that $I = \emptyset$.

Recalling that, in a category with zero object, the product of an empty family of objects is the zero object, we obtain:

e) If A is a category with products and zero object and P is a presheaf over a subspace M of C, then

P monopresheaf \longrightarrow P(θ) = 0

In particular,

 $P \text{ sheaf} \longrightarrow P(\theta) = 0$

(In the sense of {1}, a monopresheaf is a presheaf satisfying (+)). Now, we have the necessary technique in order to prove one of the crucial results of this paper.

THEOREM 2.3. If A is an A1 category and M is a subspace of a $tc^{1}C$, then for any sheaf F over X we have:

i) The sheaf F^{M} defined by $F^{M} = f^{s}f_{s}(F)$, where $f: T_{\chi} \longrightarrow T_{M}$ is the morphism ρ_{M} , satisfies $F^{M}/M \approx F/M$.

If A is an A2 category, C is, in addition, tc^2 and has a complement operator C, Cat T_X has fibered products, and X is final in Cat T_X , then for any sheaf F over X we have:

ii) If M is a closed subspace of C, then $F^M/cM \approx 0$; if the sheaf F_{cM} is defined by the exactness of the sequence $0 \longrightarrow F_{cM} \longrightarrow F \xrightarrow{\Delta_F} F^M$ (i.e. $F_{cM} = Ker \Delta_F$), then $F_{cM}/cM \approx F/cM$ and $F_{cM}/M \approx 0$.

ii) We begin with the first statement. Since cM is an open subspace of C, applying 1.17, iii it is enough to show that $F^{M}(V) = 0$, for any V open in cM. Recalling that f^{P} preserves sheaves, because f is a morphism of topologies, we see that $F^{M} \approx f^{P}i_{M}f_{s}(F)$. Therefore, we have $F^{M}(V) \approx f_{s}(F)(f(V))$; but f(V) = 0, because V is an object of T_{cM} and ρ_{M} is a representative functor (see axiom ctc3), and $f_{s}(F)(0) = 0$, because $f_{s}(F)$ is a sheaf (cf. 2.2, iii, e).Hence, $F^{M}/cM \approx 0$.

Now, we prove the second statement. Since Cat $T_{\rm X}$ has fibered products, X is a final object in Cat $T_{\rm X}$, A is an A2 category, and $\rho_{\rm A}$: $T_{\rm X} \longrightarrow T_{\rm A}$, where A is any subspace of C, preserves the "spaces" of the topologies and fibered products, then $\rho_{\rm A_S} \colon S_{\rm X} \longrightarrow S_{\rm A}$ is an exact functor (cf. {1} II, th. 4.14). Hence, $f_{\rm s}$ and $g_{\rm s}$ are exact functors (f = $\rho_{\rm M}$ and g = $\rho_{\rm CM}$).

In $S_{\mathbf{v}}$ we have the exact sequence

$$0 \longrightarrow F_{cM} \longrightarrow F \stackrel{\Delta_{\mathbf{F}}}{\longrightarrow} F^{\mathbf{M}}$$

Thus, the sequence

$$0 \longrightarrow g_{s}(F_{cM}) \longrightarrow g_{s}(F) \longrightarrow g_{s}(F^{M})$$

is exact, or equivalently, is exact the sequence

$$0 \longrightarrow F_{cM}/cM \longrightarrow F/cM \longrightarrow F^M/cM \simeq 0$$

Hence, $F_{cM}/cM \simeq F/cM$.

In a similar way, we obtain the exact sequence

 $0 \longrightarrow F_{aM}/M \longrightarrow F/M \longrightarrow F^{M}/M$

Since $\Delta_{\rm F}/{\rm M}$ is an isomorphism, the exactness of this sequence yields ${\rm F_{cM}}/{\rm M}~\simeq~0$.

The main purpose of the latter part of this section is to prove that, under certain restrictive conditions, $\Delta_F \colon F \longrightarrow F^M$ is an epimorphism. Until this moment, (C,c) will denote a fixed ctc.

If A is a category with products and zero object, we give the following

DEFINITION 2.4. If M is a subspace of C and F is a sheaf over X,we say that F is null outside M iff for all open subspace U of C we have:

$$U \times M = \theta \longrightarrow F(U) = 0$$

If M is any subspace of C, S(M) will denote the full subcategory of S_X defined by the sheaves null outside M. If f: $T_X \longrightarrow T_M$ is a morphism of topologies and A is an A1 category, $f_o: S(M) \longrightarrow S_M$ will denote the functor $f_c/S(M)$.

THEOREM 2.5. (A of type A1). If M is a subspace of C and f: $T_{X} \rightarrow T_{M}$ is the morphism ρ_{M} , then the functor for $S(M) \longrightarrow S_{M}$ has a right adjoint for $S_{M} \longrightarrow S(M)$.

Proof: Since f^s is right adjoint to f_s , it is enough to show that the image of f^s is a subcategory of S(M). (Then, f^o is f^s with S(M) as codomain).

Given a sheaf G over M, notice that $f^{s}(G) \simeq f^{p}i_{M}(G)$, where i_{M} is the inclusion functor of S_{M} in P_{M} , and so, we only need to show that

 $U \times_{\mathbf{x}} M = \theta \longrightarrow f^{p} \mathbf{i}_{M}(G)(U) = 0$

for any open subspace U of C. In fact, we have

$$f^{p}i_{M}(G)(U) = i_{M}(G)(f(U)) = i_{M}(G)(\theta) = G(\theta) = 0$$

(the last equality is true because G is a sheaf).

Now, we wish to obtain a theorem of equivalence between the categories S(M) and S_M . A similar result of Artin concerning to closed subschemes (cf. {1} III, Th. 2.2), guide us in the generalization process.

LEMMA 2.6. Let $f: K \longrightarrow K'$ be a functor and let V be an object of K' such that for any A ε ObK and any n ε Hom_V, (V, f(A)) there exist

 $U \in ObK$ and $m \in Hom_K(U,A)$ satisfying $V \xrightarrow{\sim} h f(U)$ and $f(m) \circ h = n$. Then, the full subcategory $I^f(V)$ of I^f_V defined by the class {(U,h); $U \in ObK$. $h \in Iso_K$, (V, f(U))} is initial in I^f_V .

Proof: Let (A,n) be any object in I_V^f ; applying the hypothesis on V to the morphism n: V \longrightarrow f(A), we can find a morphism in K m : U \longrightarrow A and an isomorphism in K' V \xrightarrow{h} f(U) such that f(m)h = n . Therefore, the diagram



is commutative and so, m: (U,h) \longrightarrow (A,n) is a morphism in I_{u}^{f} .

COROLLARY 2.7. If f: $K \longrightarrow K'$ is a full and representative functor, then $I^{f}(V)$ in I^{f}_{V} , for all object V of K'.

COROLLARY 2.8. (A is a complete category). Let $f: K \longrightarrow K'$ be a morphism of small categories such that K has fibered products and f is a full and representative functor which preserves fibered products. Then', any presheaf P in P(K,A) satisfies $f_p(P)(V) \approx \frac{1}{2} \frac{1}{p} p_y / I^f(V)^*$, V $\in ObK'$.

Proof: It is enough to notice that I_V^f satisfies the axiom L1*, because K has fibered products and f preserves fibered products.

DEFINITION 2.9. (C,c) is called i) ctc^{1} iff C is tc^{1} and for any closed subspace M the following conditions are satisfied:

a) Any covering in T_M is induced by ρ_M from a covering in T_X . b) If U and U' are open subspaces of C such that $U \times_X M = U' \times_X M$ and F is a sheaf null outside M, then F(U) = F(U').

ii) ctc^2 iff C is tc^2 and any closed subspace M satisfies:

 $U \times_X M = \theta \implies U \times_X cM = U$

for all open subspace U of C. iii) etc³ iff (C,c) is etc¹ and etc².

LEMMA 2.10. (A of type A0). If (C,c) is ctc^1 , Cat T_x has fibered

products, F is a sheaf null outside a closed subspace M of C, and f: $T_X \longrightarrow T_M$ is the morphism ρ_M , then the following statements are true, for any open subspace U of C:

i) $f_{p}i(F)(U \times_{X} M) \simeq F(U)$. ii) $f_{p}i(F)$ is a sheaf over M. iii) $F/M(U \times_{X} M) \simeq F(U)$.

Proof: i) Since Cat T_X has fibered products, applying 2.9 we see that $f_pi(F)(U \times_X M) \approx \underline{\lim} i(F)_{U \times_X M} / I^f (U \times_X M)^*$. Now, since (C,c) satisfies 2.9, i, b it is obvious that the values of the functor $i(F)_{U \times_X M}$ are all isomorphic, because any one is isomorphic to F(U). Hence, $\underline{\lim} i(F)_{U \times_X M} / I^f (U \times_X M)^* \approx F(U)$.

ii) Since (C,c) satisfies 2.9, i, a, applying i it follows easily that $f_pi(F)$ is a sheaf (one only needs to check the definition of sheaf).

iii) Because of ii we have $f_{g}(F) \simeq f_{p}i(F)$. Therefore, i yields the desired result.

Now, it is almost obvious how to prove:

THEOREM 2.11. (A of type A1). If (C,c) is ctc^1 , $Cat T_X$ has fiber ed products, M is a closed subspace of C, and f: $T_X \longrightarrow T_M$ is the morphism ρ_M , then the functor for $S(M) \longrightarrow S_M$ is an equivalence of categories, which inverse is $f^\circ: S_M \longrightarrow S(M)$.

Proof: By adjointness (see 2.5), there are natural transformations $\phi: e_{S(\overline{M})} \to f^{\circ}f_{\circ}$ and $\Psi: f_{\circ}f^{\circ} \longrightarrow e_{S_{\overline{M}}}$. It is a straightforward matter, which we leave to the reader, to check that ϕ and Ψ are functorial isomorphisms.

COROLLARY 2.12. If F and F' are sheaves null outside M, then

$$F/M \approx F'/M \implies F \approx F'$$

THEOREM 2.13. (A of type A1). If (C,c) is ete^2 , Cat T_x has fibe<u>r</u> ed products and M is a closed subspace of C, then

 $F/cM \approx 0 \implies F$ is null outside M

for any sheaf F over X.

Proof: Let U be an open subspace of C such that $U \times_X M = \theta$; taking in mind that (C,c) is ctc² and applying 1.17,iii we see that $F(U) \approx$ $\approx F(U \times_X cM) \approx F/cM(U \times_X cM) = 0.$

COROLLARY 2.14. If (C,c) is ctc^3 and F and F' are sheaves over X, then

$$F/M \simeq F'/M$$
, $F/cM \simeq 0 \simeq F'/cM \longrightarrow F \simeq F'$

Now, we can obtain the desired result:

THEOREM 2.15. (A of type A2). If (C,c) is ctc^3 , Cat T_X has fiber ed products, X is final in Cat T_X and M is a closed subspace of C, then the following statements are true, for any sheaf F over X :

i) $\Delta_{\mathbf{F}} \colon \mathbf{F} \longrightarrow \mathbf{F}^{\mathbf{M}}$ is an epimorphism. ii) $\mathbf{F}^{\mathbf{M}}$ is uniquely determined by F.

Proof: i) Recall that f_s and g_s are exact functors (see the proof of 2.3). Let C be the sheaf over X defined by the exactness of the sequence

$$\mathbf{F} \xrightarrow{\Delta \mathbf{F}} \mathbf{F}^{\mathbf{M}} \xrightarrow{\mathbf{C}} \mathbf{O}$$

Then, we have the exact sequence

 $g_{s}(F) \longrightarrow g_{s}(F^{M}) \longrightarrow g_{s}(C) \longrightarrow 0$

or equivalently

$$F/M \longrightarrow F^M/cM \longrightarrow C/cM \longrightarrow 0$$

and so, C/cM = 0.

In a similar way we obtain the exact sequence

 $F/M \longrightarrow F^M/M \longrightarrow C/M \longrightarrow 0$

Therefore, since Δ_F/M is an isomorphism, we conclude that C/M = 0. Now, 2.14 yields that C = 0.

ii) If F' is a sheaf over X satisfying F'/M = F/M and F'/cM = 0, then F'/M = F^M /M and F'/cM = 0 = F^M /cM. Hence, 2.14 yields that F' = F^M .

We end this section with a well known result on "characteristic" sheaves.

THEOREM 2.16. (A of type A2). If (C,c) is tc^1 and tc^2 , $CatT_X$ has fibered products, X is final in Cat T_X and M is a closed subspace of C, then for any sheaf G over M there exists a sheaf F over X such that F/M \approx G and F/CM \approx 0. If (C,c) is ctc^3 , then F is unique ly determined by G.

Proof: Since C is tc^1 , ./M: $S_X \longrightarrow S_M$ is a representative functor (see 1.13) and so, given a sheaf G over M we can find a sheaf H over X such that H/M \simeq G. Then, taking F = H^M, 2.3 enable us to conclude that F/M \simeq G and F/cM \simeq 0.

3. RELATIVE COHOMOLOGIES.

This section is devoted to realize an analysis of the cohomological effects of induced sheaves. Of course, the well known results exposed in the book of Godement (cf. {2} Ch. II, \$4.9, \$4.10, Th. 5.11.1) are obtained here, employing functorial methods. The compact exposition of cohomological theory presented in the Artin's seminar ({1} Ch. II) is continuosly used. Sheaves and presheaves are considered in this order.

I) COHOMOLOGY OF SHEAVES.

Let A be an A3 category and let C be a tc such that Cat T_X has fiber ed products and X is final in Cat T_X . (Notice that the hypothesis on Cat T_X yield the exactness of the restriction functors). We begin introducing the "true" cohomology.

DEFINITION 3.1. If M is a subspace of C, for each integer $n \ge 0$ we define the functor H_M^n : $(Cat T_X)^* \times S_X \longrightarrow A$ by:

$$H_M^n = H_{T_M}^n (. \times_X M, ./M)$$

THEOREM 3.2. The following statements are true:

If the functor ./M: $S_{\chi} \longrightarrow S_{M}$ carries injective sheaves into flask sheaves, then

iv) $H^n_M(U,) \simeq R^n H^o_M(U,)$.

If C has a complement operator c such that (C,c) is ctc^1 and M is a closed subspace of C, then for any sheaf G over M we have:

$$v$$
) $H^{n}(U \times M, G) \simeq H^{n}(U, \rho_{M}^{s}(G))$

and for any sheaf F over X we have:

$$vi$$
) $H^n_{\mathcal{M}}(U,F) \simeq H^n(U,F^M)$.

If C has a complement operator c such that (C,c) is ${\rm ctc}^3$ and M is a closed subspace of C, then for any sheaf F over X we have: vii) If F is null outside M, then ${\rm H}^{\rm n}_{\rm M}({\rm U},{\rm F}) \simeq {\rm H}^{\rm n}({\rm U},{\rm F})$. viii)There is a cohomological exact sequence of general term

$$H^{n}(U, F_{cM}) \longrightarrow H^{n}(U, F) \longrightarrow H^{n}(U, F^{M})$$

Proof: i) and ii) are trivial. iii) Notice that $H_{T_M}^*(U \times_X M,)$ is an exact cohomological functor and ./M: $S_X \longrightarrow S_M$ is an exact functor. iv) Since $R^n H_{T_M}^o(U \times_X M,) \cong H_{T_M}^n(U \times_X M,)$ and f^s , where $f:T_X \longrightarrow T_M$ is the morphism ρ_M , is an exact functor, which carries injectives into $H_{T_M}^o(U \times_X M,)$ - acyclics, the proposition follows easily: $R^n H_M^o(U,) = R^n(H_{T_M}^o(U \times_X M,) \circ f_s) \cong (R^n H_{T_M}^o(U \times_X M,)) \circ f_s \cong H_{T_M}^n(U \times_X M,) \circ f_s =$ $= H_M^n(U,)$.

v) We claim that f^s is an exact functor; since the diagram



where $j_M: S(M) \longrightarrow S_X$ is the inclusion functor, is commutative, it is enough to show that f° and j_M are exact functors. The exactness of f° is clear by reasons of equivalence (see 2.11), and the exact ness of j_M follows from the fact that S(M) is closed in S_X under taking kernels and cokernels, as it is easily deduced from the definitions.

The spectral theorem of Artin-Leray, applied to the morphism $f: T_\chi \to T_M$, tell us that

 $H^{p}(U, \mathbb{R}^{q}f^{s}(G)) \xrightarrow{p} H^{n}(U \times_{X}^{M}, G)$

for any sheaf G over M. Therefore, recalling that

$$q > 0 \longrightarrow R^q f^s = 0$$

(because f^s is exact), we obtain

_

$$H^{n}(U,f^{s}(G)) \simeq H^{n}(U \times_{v}M,G).$$

vi) Applying the above result, we have

$$H^{n}(U,F^{n}) = H^{n}(U,f^{s}f_{s}(F)) \simeq H^{n}(U \times_{Y}M,f_{s}(F)) = H^{n}_{M}(U,F).$$

vii) If (C,c) is ctc^3 , then for any sheaf F over X we have

F null outside M
$$\implies$$
 F \simeq F^M

In fact, from 2.13 follows that F^M is null outside M and so, since $F/M \,\simeq\, F^M/M$, 2.12 yields that $F \,\simeq\, F^M$.

Applying this result and i, we obtain

$$H^{n}(U,F) \simeq H^{n}(U,F^{M})$$

Hence, vi yields the desired result.

viii) By 2.3, ii and 2.15, i the sequence of sheaves over X

$$0 \longrightarrow F_{cM} \longrightarrow F \xrightarrow{\Delta_F} F^M \longrightarrow 0$$

is exact and so, iii yields the desired result.

REMARKS 3.3. i) Notice that the relative (read local) character of the cohomology just defined appears clearly in 3.2,i, 3.2,ii and 3.2,vii.

ii) Of course, the hypothesis on 3.2, iv can not be removed. Sufficient conditions in the classical case are well known (cf. {2}, II §3.3).

iii) Observe that the statement (notations as in 3.2, vi)

$$H^{n}_{cM}(U,F) \simeq H^{n}(U,F_{cM})$$

is not true, in general. Then, if we introduce the notation:

$$_{CM}H^{n}(U,F) = H^{n}(U,F_{CM})$$

under the hypothesis on 3.2, viii, we obtain an exact cohomological sequence of general term

$$_{cM}H^{n}(U,F) \longrightarrow H^{n}(U,F) \longrightarrow H^{n}_{M}(U,F)$$

Now, we focus our attention in the cohomology with presheaves values. DEFINITION 3.4. If M is a subspace of C, for each integer $n \ge 0$ we define the functor H_M^n : $S_X \longrightarrow P_M$ by:

 $H_{M}^{n} = H_{T_{M}}^{n}(./M)$

THEOREM 3.5. The following statements are true:

i) $F/M \simeq F'/M \Longrightarrow H^n_M(F) \simeq H^n_M(F')$.

 $ii) \quad H_{X}^{n} \simeq H_{T_{X}}^{n}$

iii) H_w is an exact cohomological functor.

If the functor ./M: $S_{\chi} \longrightarrow S_M$ carries injective sheaves into flask sheaves, then

iv) $H_{M}^{n} \simeq F^{n}H_{M}^{\circ}$

Without assumptions, we have for any sheaf F over X:

$$v$$
) $H^n_M(F)_{\circ \rho_M} \simeq H^n_M(,F)$.

If C has a complement operator c such that (C,c) is ctc^1 and M is a closed subspace of C, then for any sheaf F over X we have:

vi) $H^n_M(F) \circ \rho_M \simeq H^n(F^M)$.

If (C,c) is ctc³, we also have:

vii) If F is null outside M, then $H^n_M(F)_{\circ \rho_M} = H^n(F)$.

viii) There is an exact cohomological sequence of general term

$$H^{n}(F_{cM}) \longrightarrow H^{n}(F) \longrightarrow H^{n}(F^{M})$$

Proof: i), ii), iii) and iv) can be obtained as in 3.2. v) Knowing that $\#^n_{T_M}(G) \approx \#^n_{T_M}(G)$, for any sheaf G over M, the proposition follows easily:

 $\begin{aligned} & \mathcal{H}^{n}_{M}(F)\left(U\times_{X}M\right) = \mathcal{H}^{n}_{\mathcal{T}_{M}}(F/M)\left(U\times_{X}M\right) \simeq \mathcal{H}^{n}_{\mathcal{T}_{M}}\left(U\times_{X}M,F/M\right) = \mathcal{H}^{n}_{M}(U,F). \\ & \text{vi) Applying v and 3.2,vi, we obtain:} \\ & \mathcal{H}^{n}_{W}(F)\left(U\times_{v}M\right) \simeq \mathcal{H}^{n}_{W}(U,F) \simeq \mathcal{H}^{n}\left(U,F^{M}\right) \simeq \mathcal{H}^{n}\left(F^{M}\right)\left(U\right). \end{aligned}$

vii) The statement in question can be obtained as vi, applying now v and 3.2, vii. Also, it can be proved in the following way: since (C,c) is ctc³, for any sheaf F over X we have

F null outside M
$$\longrightarrow$$
 F \simeq F^M

Hence, $H^{n}(F) \simeq H^{n}(F^{M})$ and so, the proposition follows from vi. viii) It can be obtained as 3.2, vii.

II) COHOMOLOGY OF PRESHEAVES.

Let A be an A1 category and let C be an arbitrary tc. First, we consider the cohomology of a covering. In order to conserve a spectral result and to obtain a new one, we adopte the following

DEFINITION 3.6. If M is a subspace of C and K_U is the category of coverings of an open subspace U of C, for each integer $n \ge 0$ we define the functor $\Pi_M^n: K_U^* \times P_M \longrightarrow A$ by

$$H_{M}^{n} = H_{\mathcal{T}_{M}}^{n} (\cdot \times_{X}^{M}, \cdot)$$

THEOREM 3.7. The following statements are true:

 $\begin{array}{ll} i) & (U_{i} \rightarrow U)_{i \in I} \cong (U_{j}' \rightarrow U)_{j \in J} \implies H^{n}_{M}((U_{i} \rightarrow U)_{i \in I},) \cong H^{n}_{M}((U_{j}' \rightarrow U)_{j \in J},) \\ ii) & H^{n}_{X} \cong H^{n}_{T_{X}} \\ iii) & H^{*}_{M}((U_{i} \rightarrow U)_{i \in I},) \text{ is an exact cohomological functor.} \\ If A is an A3 category, then we have: \\ \end{array}$

$$iv \to H^{n}_{M}((U_{i} \longrightarrow U)_{i \in I},) \approx R^{n}H^{o}_{M}((U_{i} \longrightarrow U)_{i \in I},)$$

$$v \to H^{p}_{M}((U_{i} \longrightarrow U)_{i \in I}, H^{q}_{M}(F)) \xrightarrow{p} H^{n}_{M}(U,F).$$

If A is an Al category, C has a complement operator c such that (C,c) ctc^{1} , Cat T_{X} has fibered products and M is a closed subspace of C, then for any sheaf F over X we have:

vi) If F is null outside M, then
$$H^n_M((U_i \longrightarrow U)_{i \in I}, F/M) \simeq H^n_M((U_i \longrightarrow U)_{i \in I}, F)$$

If A is an A3 category and X is final in Cat T_X , we also have:
vii) If F is null outside M, then

$$H^{p}_{M}((U_{i} \longrightarrow U)_{i \in I}, \rho_{M_{p}}(H^{q}(F))) \xrightarrow{p} H^{n}(U, F)$$

Proof. i) and ii) are trivial.

iii) Notice that $H^*_{T_M}((U_i \times_X M \longrightarrow U \times_X M)_{i \in I},)$ is an exact cohomolo-

gical functor.

iv) Since A is an A3 category, $H^n_{T_M}((U_i \times_X^M \longrightarrow U \times_X^M)_{i \in I},) \approx \mathbb{R}^n H^o_{T_M}((U_i \times_X^M \longrightarrow U \times_X^M)_{i \in I},).$

v) Since A is an A3 category, the cohomologies of sheaves are $d\underline{e}$ fined and we have

$$H^{\mathbf{p}}_{\mathcal{T}_{\mathbf{M}}}((U_{\mathbf{i}} \times_{\mathbf{X}}^{\mathbf{M}} \longrightarrow U \times_{\mathbf{X}}^{\mathbf{M}})_{\mathbf{i} \in \mathbf{I}}, H^{\mathbf{q}}_{\mathcal{T}_{\mathbf{M}}}(F/M)) \longrightarrow H^{\mathbf{n}}_{\mathcal{T}_{\mathbf{M}}}(U \times_{\mathbf{X}}^{\mathbf{M}}, F/M)$$

vi) It follows easily from 2.10, iii, by a direct analysis of the complex which gives the cohomology.

vii) Observe that S(M) has injectives, because it is equivalent to the category of sheaves over M (see 2.11) (since A is an A3 catego ry, ${}^{S}_{M}$ has injectives). Also, observe that the functor j_{M} carries injectives into flasks, because $f^{S} = j_{M} \circ f^{\circ}$, where f: $T_{X} \longrightarrow T_{M}$ is the morphism ρ_{M} , and f^{S} has this property.

Now, consider the (two) functors given by the commutative diagram



Let us evaluate its derived functors. By iv, we have

$$\mathbb{R}^{P} \mathbb{H}_{M}^{\circ}((\mathbb{U}_{i} \longrightarrow \mathbb{U})_{i \in \mathbb{I}},) \simeq \mathbb{H}_{M}^{P}((\mathbb{U}_{i} \longrightarrow \mathbb{U})_{i \in \mathbb{I}}, \mathbb{U})$$

Recalling that f_p is an exact functor, by the hypothesis on Cat τ_{χ} , and that j_M is an exact functor, which carries injectives into i-acy clics, we obtain

$$R^{q}(f_{p}ij_{M}) \approx f_{p} \circ R^{q}(ij_{M}) \approx f_{p} \circ (R^{q}i) \circ j_{M} = f_{p} \circ H^{q} \circ j_{M}$$

Both results elucidate the first member of the spectral convergence in question. Concerning to the second member, 2.10,i implies

$$\begin{split} H^{\circ}_{M}((U_{i} \longrightarrow U)_{i \in I} ,) \circ (f_{p}ij_{M}) &\simeq H^{\circ}((U_{i} \longrightarrow U)_{i \in I} ,) \circ (ij_{M}) \\ \text{Therefore, recalling that } H^{\circ}((U_{i} \longrightarrow U)_{i \in I} ,) \circ i &\simeq \Gamma_{U} , \text{ we obtain} \end{split}$$

$$H_{M}^{\circ}((U_{i} \rightarrow U)_{i \in I},) \circ (f_{p}ij_{M}) \simeq \Gamma_{U} \circ j_{M}$$

and so, since j_M is an exact functor which carries injectives into Γ_{II} -acyclics, we have

$$R^{n}(H^{\circ}_{M}((U_{i} \rightarrow U)_{i \in I},) \circ (f_{p}ij_{M})) \approx R^{n}(\Gamma_{U}j_{M}) \approx (R^{n}\Gamma_{M}) \circ j_{M} = H^{n}(U,) \circ j_{M}$$

We introduce the limit cohomology of presheaves by a more general procedure than the one used by Artin in $\{1\}$. Of course, both definitions agree in the case that the category of values is A3.

Let A be an A1 category and let T be an arbitrary topology. If U is an object of Cat T and K is a subcategory of K_U , for each integer $n \ge 0$ we define the functor $H^n_T(K,): P(T,A) \longrightarrow A$ by:

$$H_T^{n}(,P) = \underline{\lim} H_T^{n}(,P) \circ k^*, P \in ObP$$

where k: $K \longrightarrow K_U$ is the inclusion functor (notice that $H^n_T(, P)$: $K^*_U \longrightarrow A$). It is straightforward, to check the following propositions:

i) If A is an A2 category and K is filtrant, then $H_T^*(K,)$ is an exact cohomological functor.

ii) If A is an A3 category and K is filtrant, then $H^n_T(K,) \simeq R^n H^o_T(K,)$. (Concerning to i, the usual statement about the exactness of the limit is required; and for ii, the proposition i tell us that it is enough to show that $H^n_T(K,)$ vanishes on injectives, if n > 0).

Notice that all the other results of {1}, concerning to limit cohomo logy, are preserved by our definition.

The Cech cohomology of presheaves is introduced following {1}. Now, we focus our attention in the relative limit cohomology.

DEFINITION 3.8. If M is a subspace of C and K is a subcategory of K_U for each integer $n \ge 0$ we define the functor $H^n_M(K,): P_M \longrightarrow A$ by

$$H^{n}_{M}(K,) = H^{n}_{T_{M}}(K \times M,)$$

(Notice that K_{X}^{M} is a subcategory of $K_{U}^{X}{}_{X}^{M}$, which is a subcategory of $K_{U}{}_{X}{}_{M}$).

THEOREM 3.9. The following statements are true: i) $K \times_X M = K' \times_X M \longrightarrow H^n_M(K,) = H^n_M(K',).$ ii) $H^n_X(K,) = H^n_{T_X}(K,).$ If k: $K \longrightarrow K_U$ is the inclusion functor, then for any presheaf P aver M we have iii) $H^n_M(K,P) = \lim_{M \to M} H^n_M(,P) \circ k^*$ If K^* is filtrant, then: iv) $H^o_M(K,) \circ i_M = \Gamma_{U \times_V M}$ If A is an A2 category and K^* is filtrant, then: v) $H_M^*(K,)$ is an exact cohomological functor. If A is an A3 category and K^* is filtrant, then: vi) $H_M^n(K,) \simeq R^n H_M^o(K,)$. vii) $H_M^p(K, H_M^q(F)) \xrightarrow{P} H_M^n(U, F)$

If A is an Al category, K^* is filtrant, C has a complement operator C such that (C,C) is cto^1 , Cat T_X has fibered products and M is a closed subspace of C, then for any sheaf F over X we have: viii)If F is null outside M, then $\operatorname{H}^n_M(K,F/M) = \operatorname{H}^n(K,F)$. If A is an A3 category and X is final in Cat T_X , then we also have: ix) If F is null outside M, then $\operatorname{H}^p_M(K,\rho_M \operatorname{H}^q(F)) \longrightarrow \operatorname{H}^n(U,F)$.

Proof: i) and ii) are trivial. iii) If $k_M : K \times_X M \longrightarrow K_{U \times_X M}$ is the inclusion functor, by definition we have for any presheaf P over M

$$H^{n}_{\mathcal{T}_{M}}(K \times_{X}^{M}, P) \simeq \underline{\lim} H^{n}_{\mathcal{T}_{M}}(, P) \circ k_{M}^{*}$$

and it is clear that

$$\underbrace{\lim_{T_{M}} H^{n}_{T_{M}}}_{I,M}(,P) \circ k^{*}_{M} \simeq \underbrace{\lim_{T_{M}} H^{n}_{T_{M}}}_{M}(.\times_{X}^{M},P) \circ k^{*}$$

iv) Notice that

K* filtrant ----> K×_xM* filtrant

and so, we have

 $H^{\circ}_{T_{M}}(K \times_{X} M,) \circ i_{M} \simeq \Gamma_{U \times_{X} M}$

v) Since $K \times_X M^*$ is filtrant, then $H_{T_M}^*(K \times_X M,)$ is an exact cohomological functor.

vi) Recalling that K* is filtrant, we have

$$H^{n}_{\mathcal{T}_{M}}(K \times_{X} M,) \simeq R^{n} H^{\circ}_{\mathcal{T}_{M}}(K \times_{X} M,)$$

vii) By the same reasons, we have the spectral convergence

$$H^{p}_{\mathcal{T}_{M}}(K \times_{X}^{M}, \mathcal{H}^{q}_{\mathcal{T}_{M}}(F/M)) \xrightarrow{p} H^{n}_{\mathcal{T}_{M}}(U \times_{X}^{M}, F/M).$$

viii) Applie 3.7, vi and pass to the limit over K_U^* , using the proposition iii.

ix) It can be obtained as 3.7, vii.

We end this section introducing the relative Cech cohomology of pre

sheaves. The definition is not the expected one, because the nat<u>u</u> ral definition do not preserves the relative character (see 3.11,i). However, in the special case 3.11,xi both procedures agree.

DEFINITION 3.10. If M is a subspace of C and U is open, for each integer $n \ge 0$ we define the functor $\check{H}^n_M(U,): P_M \longrightarrow A$ by:

$$\check{H}^{n}_{M}(U,) = \check{H}^{n}_{T_{M}}(U \times_{X} M,)$$

THEOREM 3.11. The following statements are true: $\mathbb{U}_{X}^{M} \simeq \mathbb{U}_{X}^{M} \longrightarrow \tilde{\mathbb{H}}_{M}^{n}(\mathbb{U},) \simeq \tilde{\mathbb{H}}_{M}^{n}(\mathbb{U},)$ i) $\tilde{H}_{X}^{n}(U,) \approx \tilde{H}_{T_{X}}^{n}(U,)$ ii) *iii)* $\tilde{H}^{n}_{M}(U,P) \approx \underline{\lim}_{M} H^{n}_{M}(,P)$ *iv)* $\tilde{H}_{M}^{\circ}(U,) \circ i_{M} \simeq \Gamma_{U \times xM}$ If A is an A2 category, then: $\check{H}^{\star}_{M}(U,)$ is an exact cohomological functor. v) If A is an A3 category, then: vi) $\check{H}^{n}_{M}(U,) \approx R^{n}\check{H}^{o}_{M}(U,)$ $\begin{array}{l} \textit{viii} \quad \overset{}{H}^{p}_{M}(U, H^{q}_{M}(F)) \xrightarrow{P} H^{n}_{M}(U, F) \\ \textit{viii} \quad \overset{}{H}^{1}_{M}(U, F/M) \simeq H^{1}_{M}(U, F) \quad , \quad \overset{}{H}^{2}_{M}(U, F/M) \subset H^{2}_{M}(U, F) \end{array}$ If A is an Al category, C has a complement operator c such that (C,c) is ctc^{1} , Cat T_x has fibered products and M is a closed subspace of C, then for any sheaf F over X we have: If F is null outside M, then $\check{H}^{n}_{M}(U,F/M) \simeq \check{H}^{n}(U,F)$ ix) If A is an A2 category and X is final in Cat T_x , then we also have: If F is null outside M, then $\check{H}^{P}_{M}(U,\rho_{M_{n}}(H^{q}(F))) \xrightarrow{p} H^{n}(U,F)$. x) If A is an Al category and M is a subspace of C such that any cove \underline{r} ing of $U \star_X^M$ is induced by ρ_M from a covering of U, then we have: xi) $H^{n}_{M}(U,) \simeq H^{n}_{M}(K_{u},).$

Proof: i), ii), iii), iv), v), vi) and vii) can be obtained as the homologous propositions of 3.9.

viii) It follows immediately from

$$\check{\mathrm{H}}_{T_{\mathrm{M}}}^{1}(\mathrm{U}_{\mathrm{X}}^{\mathrm{M}},\mathrm{F/M}) \simeq \mathrm{H}_{T_{\mathrm{M}}}^{1}(\mathrm{U}_{\mathrm{X}}^{\mathrm{M}},\mathrm{F/M}) \ , \quad \check{\mathrm{H}}_{T_{\mathrm{M}}}^{2}(\mathrm{U}_{\mathrm{X}}^{\mathrm{M}},\mathrm{F/M}) \subset_{\sim} \mathrm{H}_{T_{\mathrm{M}}}^{2}(\mathrm{U}_{\mathrm{X}}^{\mathrm{M}},\mathrm{F/M}) \ .$$

It should be pointed out that viii could be obtained from vii just as in the absolute cohomology case. xi) The hypothesis on M tell us that $K_U \times_X M \simeq K_{U \times_V M}$. Hence, we have

$$H^{n}_{\mathcal{T}_{M}}(K_{U\times_{X}M},) \simeq H^{n}_{\mathcal{T}_{M}}(K_{U}\times_{X}M,)$$

or equivalently

$$\tilde{H}_{M}^{n}(U,) \simeq H_{M}^{n}(K_{U},)$$

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REVISTA DE LA UNION MATEMATICA ARGENTINA Volumen 24, Número 2, 1968

NOTE ON GALOIS EXTENSION OVER THE CENTER by Manabu Harada

In {2} S.U.Chase, D.K.Harrison and A.Rosenberg obtained a Galois Theory for strongly Galois extensions of commutative rings (CHR-Galois). This was generalized to non commutative rings by F.R. Demeyer {3} , T.Kanzaki {6} , H.F.Kreimer {8} and others. Recent ly, O.E.Villamayor and D.Zelinsky obtained in {11} a weak Galois theory of commutative rings in order to study the strong one from a different point of view.

In the first section of this short paper we shall use similar arguments to those of {11} to show that if an algebra Λ over a commutative ring R is a strongly Galois extension of R, then Λ and its center C are weakly Galois extensions over C and R, respectively. If Λ is a weakly Galois extension over C, Λ is the sum of all C-modules J_{σ} (see below or {10} for the definition of J_{σ}). By means of some properties of the J_{σ} 's, we shall study in section 2, a Galois theory over the center, the argument being similar to that of {7} and {5}, Theorem 1.

The author would like to express his thanks to Professor O.E.Vil<u>la</u> mayor for inviting him to Universidad de Buenos Aires and giving an opportunity to see his and Zelinsky's preprint of {11}.

1. GALOIS EXTENSION OVER R.

Let R be a commutative ring with identity, Λ an algebra over R, C the center of Λ and G a finite group of automorphisms of Λ . We say that Λ is a Galois extension with respect to G of its G-fixed ring Λ^{G} if there exist elements x_{i} , y_{i} i=1,2,...,n, in Λ such that $\sum_{\sigma \in G} \sigma(x_{i})y_{i} = \delta_{\sigma,1}$. We note that if $\Lambda^{G} = R$, then Λ is a finitely generated and separable R-algebra by {9}, Lemma 2.

Let Γ be an R-subalgebra of Λ and G the group of all automorphisms of Λ leaving invariant the elements of Γ . We quote here the definition of weakly Galois extension of {11}: Λ is said to be a (right) weakly Galois extension of Γ if the following two conditions are satisfied:

a) A is a finitely generated projective right r-module.

b) $G\Lambda_{\ell} = Hom_{\Gamma}(\Lambda, \Lambda)$, where $\sum g_{i}x_{i\ell}(\lambda) = \sum g_{i}(x_{i\lambda})$ for $g_{i} \in G, x_{i,\lambda} \in \Lambda$.

LEMMA 1. Let Λ be a strongly Galois extension of R with group G. If R has no proper idempotents, then Λ and its center C are weakly

Galois extensions of C and R, respectively.

Proof: Since A is a finitely generated R-module, there exist mutual ly orthogonal primitive central idempotents e_i such that $\sum_{i=1}^{n} e_i = 1$. Then $\operatorname{Hom}_{C}(\Lambda,\Lambda) = \sum_{i} \oplus \operatorname{Hom}_{Ce_i}(\Lambda e_i,\Lambda e_i)$ and $\operatorname{Hom}_{R}(C,C) = \sum_{i} \oplus \operatorname{Hom}_{Re_i}(Ce_i,Ce_i)$. Let $T_i = \{g \mid cG, g(e_i) = e_i\}$. We can easily see that Λe_i is a strongly Galois extension of Re_i with group T_i (see $\{3\}$). Furthermore, Λe_i and Ce_i are strongly Galois extensions over Ce_i and Re_i with respect to H_i and T_i/H_i by $\{3\}$ and $\{9\}$, where $H_i = \{g \mid cT_i, g(d) = c \lor ccCe_i\}$. Hence $\operatorname{Hom}_{Ce_i}(\Lambda e_i, \Lambda e_i) = H_i(\Lambda e_i)_{\ell}$ and $\operatorname{Hom}_{Re_i}(\operatorname{Ce}_i, \operatorname{Ce}_i) = (T_i/H_i)(\operatorname{Ce}_i)_{\ell}$. Now, we put $H_i^* = \{h \mid cG, h \mid \Lambda e_i = h'$ for some $h' cH_i, h \mid \Lambda e_j = I_{\Lambda e_j}$ for $i \neq j$. Then $H_i^* \subseteq H = \{h \mid cG, h \mid C = I_C\}$. Therefore, $\operatorname{Hom}_C(\Lambda,\Lambda) = H\Lambda_{\ell}$. Similarly we obtain $\operatorname{Hom}_R(C,C) = G'C_{\ell}$, where G' is the group of all automorphisms of C over R. Condition a) follows from $\{1\}$, since Λ is separable over R.

THEOREM 1. Let Λ be a strongly Galois extension of R with finite group G. Then Λ and its center C are weakly Galois extensions of C and R, respectively.

Proof: We shall use the same notation and argument of {11}. Since Λ is $\Lambda \otimes \Lambda^*$ -projective, $C_x = (\operatorname{Hom}_{\Lambda}{}^{\ell}(\Lambda,\Lambda))_x = \operatorname{Hom}_{\Lambda}{}^{\ell}(\Lambda_{\alpha},\Lambda_{\alpha})$ (cf.{11}, (2.7)). Furthermore, R_x has no proper idempotents by {11}, (2.13). Hence, $H(x)(\Lambda_x)_{\ell} = \operatorname{Hom}_{C_x}(\Lambda_x,\Lambda_x)$ and $G'(x)(C_x)_{\ell} = \operatorname{Hom}_{R_x}(C_x,C_x)$, where H(x), G'(x) are as above in Λ_x and C_x . Since C is R-finitely generated, all elements of H(x) and G'(x) are induced by elements of H and G', respectively (by {11}, (2.14)). Hence $(H\Lambda\ell)_x =$ $= H(x)(\Lambda_x)_{\ell} = \operatorname{Hom}_{C_x}(\Lambda_x,\Lambda_x) = (\operatorname{Hom}_C(\Lambda,\Lambda))_x$ since Λ is C-finitely generated and projective. Therefore, $H\Lambda_{\ell} = \operatorname{Hom}_C(\Lambda,\Lambda)$. Similarly $G'C_{\ell} = \operatorname{Hom}_{P}(C,C)$.

We shall give latter an example in which a strongly Galois extension of R is not a strongly Galois extension over its center with respect to the corresponding subgroup.

2. GALOIS EXTENSION OVER CENTER.

In this section we always assume that Λ is a separable algebra over its center C. If $H\Lambda_{\ell} = Hom_{C}(\Lambda,\Lambda)$, then for any element x in Λ^{H} , x_{ℓ} belongs to the center of $Hom_{C}(\Lambda,\Lambda) = C$; hence $\Lambda^{H} = C$. We shall study some properties are treated in {7} and {10}. For $\sigma \in H$ let $J_{\sigma} =$ $= \{x \mid \epsilon \Lambda$, $yx = x\sigma(y)$ for all $y\epsilon\Lambda\} = Hom_{\Lambda}\ell(\Lambda,\Lambda\sigma)$, where $\Lambda\sigma$ is the same module as Λ as left Λ -module and the operation of Λ as right Λ-module is defined by x*y = xσ(y). Furthermore, Λ_σ = Λ Θ_CJ_σ and Λ = ΛJ_σ (see {10}). For any element $x_t \sigma \in J_{\sigma t} \sigma \subset$ Λ_tH(x_tσ)(y) = x_{σ(y)} = yx_r for every y ∈ Λ. Hence, (J_σ)^σ_t = (J_σ)_r and Λ_tσ = Λ_tJ_{σt}σ = Λ_t(J_σ)_r = Λ_t Θ_C(J_σ)_r since Λ_tΛ_r = Λ_t Θ_CΛ_r by {1}.

PROPOSITION 2. Let Λ be separable over its center C and S a subset of H. Then $S\Lambda_{\ell} = Hom_{C}(\Lambda, \Lambda)$ if and only if $\Lambda = \sum_{\sigma \in S} J_{\sigma}$.

Proof: $S\Lambda_{\ell} = \sum_{\sigma \in S} \Lambda_{\ell} \sigma = \sum_{\sigma} \Lambda_{\ell} \otimes (J_{\sigma})_{r}$. Since C is a C-direct summand of Λ and $Hom_{C}(\Lambda,\Lambda) = \Lambda_{\ell} \otimes_{C} \Lambda_{r}$, the proposition follows.

COROLLARY. A is a weakly Galois extension if and only if $\Lambda = \sum_{\sigma \in S} J_{\sigma}$, where S is a finite subset of H. Furthermore, A is generated by units as C-module if and only if $SA_{\ell} = Hom_{C}(\Lambda, \Lambda)$ and the elements of S are inner-automorphisms.

Proof: It is clear.

Let S be a subset of H. We call S strongly distinct if there exists a family of elements $\{x_i^{(\sigma)}, y_i^{(\sigma)}\}_{i=1}^{n=n(\sigma)}, \sigma \in S$ such that

 $\sum_{i} \tau(x_{i}^{(\sigma)}) y_{i}^{(\sigma)} = \delta_{\tau,\sigma} \text{ for all } \sigma, \tau \in S.$

It is clear that this condition is equivalent with the existence of Galois generators if S is a group.

THEROREM 3. Let S, J_{σ} be as above and $\Gamma = \sum_{\sigma \in S} J_{\sigma}$. If S is strongly distinct, then $\Gamma = \sum_{\sigma \in S} \Theta J_{\sigma}$. Conversely, if $\Gamma = \sum_{\sigma \in S} \Theta J_{\sigma}$ and Γ is a direct summand of Λ as C-module, then S is strongly distinct.

Proof: Assume that $\Gamma = \sum_{\sigma} \oplus J_{\sigma}$ and Γ is a direct summand of Λ as C-module. $\Lambda_{\ell} \otimes \Lambda_{\Gamma} = \operatorname{Hom}_{\mathbb{C}}(\Lambda, \Lambda)$, since Λ is C-separable. Let P_{σ} be a projection of Λ onto J_{σ} . Then $P_{\sigma} \in \operatorname{Hom}_{\mathbb{C}}(\Lambda, \Lambda)$. Hence, there exist elements $\{x_{i}^{(\sigma)}, y_{i}^{(\sigma)}\}_{i=1}^{n=n(\sigma)}$ in Λ such that $\sum_{i\ell} x_{i\ell}^{(\sigma)} \otimes y_{i\Gamma}^{(\sigma)} = P_{\alpha}$. There fore, $0 = P_{\sigma}(J_{\tau}) = J_{\tau}\sum_{\tau} \tau(x_{i}^{(\sigma)})y_{i}^{(\sigma)}$ for $\sigma \neq \tau$. Since $\Lambda J_{\tau} = \Lambda$, $\sum_{\tau} \tau(x_{i}^{(\sigma)})y_{i}^{(\sigma)} = 0$. Similarly, $J_{\sigma}(I - \sum_{\sigma} \sigma(x_{i}^{(\sigma)})y_{i}^{(\sigma)}) = 0$. Hence, $I = \sum_{\sigma} \sigma(x_{i}^{(\sigma)})y_{i}^{(\sigma)}$. Conversely, assume S is strongly distinct. We assume $0 = \sum_{\sigma \in S} z_{\sigma}$, $z_{\sigma} \in J_{\sigma}$. Then $0 = \sum_{i}\sum_{\tau} \tau x_{i}^{(\sigma)} z_{\tau} y_{i}^{(\sigma)} = \sum_{\sigma} z_{\sigma}(\sum_{i} \tau(x_{i}^{(\sigma)})y_{i}^{(\sigma)}) = z_{\sigma}$. Hence, $\Gamma = \sum_{\sigma} \oplus J_{\sigma}$.

LEMMA 2. Let $A \supset B$ be R-algebras. If B is R-separable, then $V_A(B)$

is a direct summand of A as two-sided $V_A(B)$ -module, where $V_A(B) = {a | \epsilon A , ba = ab for all b \epsilon B}$.

Proof: We consider A as a left B $\Theta_{R}B^{*}$ -module. Since B is R-separable, there exist elements $\{x_{i}, y_{i}\}_{i}$ in A such that $[x_{i}y_{i}] = 1$ and $[bx_{i} \Theta y_{i}] = [x_{i} \Theta y_{i}] b$ for all $b \in B$. Now we define a map $\phi: A \to A$ by setting $\phi(a) = [(x_{i} \Theta y_{i}^{*})a] = [x_{i}ay_{i}]$ for a ϵ A. From the above relations of $\{x_{i}, y_{i}\}$ we obtain $\phi(A) \subseteq V_{A}(B)$ and $\phi|V_{A}(B)] = I_{V_{A}(B)}$. Furthermore, $\phi(aa') = [x_{i}aa'y_{i}] = ([x_{i}ay_{i})a'] = \phi(a)a'$ for $a' \epsilon V_{A}(B)$. Similarly $\phi(a'a) = a'\phi(a)$. Hence, ϕ is $V_{A}(B) - V_{A}(B)$ homomorphism. Since $V_{A}(B)$ is $V_{A}(B)$ - projective, $A = V_{A}(B) \oplus \ker \phi$.

PROPOSITION 4. Let Λ be a central separable C-algebra and Γ a separable subalgebra ($\Gamma \supseteq C$). Then Γ is a direct summand of Λ as a two-sided Γ -module.

Proof: We know from {6}, Theorem 2 that $r = V_{\Lambda}$ (V_{Λ} (r)) and V_{Λ} (r) is C-separable. Hence the proposition follows from Lemma 2. From now we assume that the subset S of H is a finite group G.

PROPOSITION 5. Let Λ be a central separable C-algebra and G a finite subgroup of the group of C-automorphism of Λ_3 let $\Gamma = \sum_{\sigma \in G} J_{\sigma}$. Then the following statements are equivalent:

1) $\Gamma = \sum_{\sigma \in G} \mathfrak{G}_{\sigma} \text{ and } |G| \text{ is a unit in C.}$ 2) $\Gamma = \sum_{\sigma \in G} \mathfrak{G}_{\sigma} \text{ and } \Gamma \text{ is C-separable.}$ 3) $\Lambda \text{ is a strongly Galois extension of } \Lambda^{G} \text{ and } \Lambda^{G} \text{ is C-separable.}$

Proof: 1) \leftrightarrow 2) It is clear from ({4} Lemma 4) by localization of C.

2) \longrightarrow 3) Since r is C-separable, r is a direct summand of A as r-module by Proposition 4. Hence, G is strongly distinct by Theorem 3.

3) \longrightarrow 1) |G| is unit in C by {5} , Proposition 5, and the rest is clear.

LEMMA 3. $J_{\sigma}J_{\tau} = J_{\tau\sigma}$ for any σ, τ .

Proof: Let m be a maximal ideal in C. Then $(J_{\sigma})_{m} = \text{Hom}_{\Lambda_{m}^{\mathfrak{l}}}(\Lambda_{m}, \Lambda_{m}^{\sigma})^{\mathfrak{l}} = C_{m}u_{\sigma}$, where $\sigma(y) = u_{\sigma}^{-1}yu_{\sigma}$. Hence, $(J_{\sigma}J_{\tau})_{m} = C_{m}u_{\sigma}u_{\tau}$ and $u_{\tau\sigma}(u_{\sigma}u_{\tau})^{-1}$,

 $u_{\sigma}u_{\tau}u_{\tau\sigma}^{-1}$ belong to C_{m} . Therefore, $J_{\sigma}J_{\tau} = J_{\tau\sigma}$.

PROPOSITION 6. Let G be a finite subgroup of H and assume Λ is a strongly Galois extension of Λ^{G} ; let Ω be a separable C-algebra b<u>e</u> tween Λ and Λ^{G} . Then the following statements are equivalent: 1) $\Omega = \Lambda^{H}$ for some subgroup H of G. 2) $V_{\Lambda}(\Omega) = \sum_{\sigma \in S} J_{\sigma}$ for some subset S of G. 3) There exist elements $\{x_i \in \Omega, y_i \in \Lambda\}$ such that $[x_iy_i = I and$ $\sum \rho(\mathbf{x}_i) \mathbf{y}_i = 0 \text{ for } \rho | \Omega \neq \mathbf{I}_{\Omega}, \ \rho \in \mathbf{G}. \quad (\{8\}, \text{ Proposition } 3.5).$ *Proof:* Since $V_{\Lambda}(\Omega) \subseteq \sum_{\sigma \in G} \Phi J_{\sigma}$, S is a subgroup from Lemma 3. Fur thermore, $\operatorname{Hom}_{\Omega_{-}}(\Lambda,\Lambda) = \Lambda_{\ell} \Theta_{C} V_{\Lambda}(\Omega)_{r}$ by {6}, Theorem 2, since Ω is C-separable and therefore 1) and 2) are equivalent. 1) \longrightarrow 3) Let G = $\bigcup_{i} \rho_{i}H$. Then $\Gamma = \sum_{\sigma \in G} \Theta J_{\sigma} = \sum_{i} \Gamma_{H} J_{\rho_{i}}$ and Γ is a direct summand of Λ as r-module, where $r_{\rm H} = \sum_{\sigma \in \rm H} \Theta J_{\sigma}$. Let p be a projection of Λ onto $\Gamma_H J_i = J_H$. Then $p \in \operatorname{Hom}_{\Gamma H_g}(\Lambda, \Lambda) = (\Lambda^H)_g \otimes \Lambda_r$. Hence, there exist { $x_i \in \Lambda^H$, $y_i \in \Lambda$ } such that $\sum x_{i\ell} \otimes y_{ir} = p$. 3) \longrightarrow 2) Put H = { $\sigma | \epsilon G, \sigma | \Omega = I_{\Omega}$ }. Then $V_{\Lambda} (\Omega) \supseteq \sum_{\sigma \epsilon H} J_{\sigma}$. Let $y \in V_{\Lambda}$ (Ω) and $y = x_{\rho_1} + x_{\rho_2} + \dots$, where $x_{\rho_i} \in \Gamma_H J_{\rho_i}$. Then $y = yI = V_{\Lambda}$ $= y \sum_{i} x_{i} y_{i} = \sum_{i} x_{i} y_{i} = \sum_{j} \sum_{i} x_{i} x_{\rho_{j}} y_{i} = \sum_{\rho_{j}} \sum_{\rho_{j}} (x_{i}) y_{i} = x_{\rho_{1}}.$ Hence $V_{\Lambda}(\Omega) =$ $= \sum_{\alpha \in H} J_{\alpha}$.

Finally, we shall give an example of a strongly Galois extension Λ of R, such that Λ is not a strongly Galois extension over its center with respect to its subgroup. However Λ is a strongly Galois extension over its center with respect to a suitable group.

Let G_2 be a cyclic group of order 2 and Q the field of rational numbers. Put $G = G_2 \times G_2$ and $K = Q(\sqrt{2})$. Then $L = K \otimes_Q K$ is a strongly Galois extension of Q with respect to G by {9}, Proposition 1. Let g and h be the inner-automorphisms of Q₂ induced by

 $\begin{pmatrix} -1 & 0 \\ & \\ 0 & 1 \end{pmatrix} \quad \text{and} \quad \begin{pmatrix} 0 & 1 \\ & \\ 1 & 0 \end{pmatrix}$

respectively. Then (g) × (h) \approx G and Q₂ is a strongly Galois extension of Q with respect to G, since $\frac{1}{2}$ { e_{11} , e_{11} , e_{22} , e_{22} , e_{21} , e_{12} , e_{12} , e_{21} } is a family of Galois generators. Put $\Lambda = Q_2 \oplus L$. Then Λ is a strong ly Galois extension of Q with group G if we define g(a+b) = g(a)+g(b)

for $g \in G$, $a \in Q$ and $b \in L$. It is clear that the fixing group of its center is equal to G. But if we define g(a + b) = g(a) + b, then A is a strongly Galois extension of its center with respect to G.

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A MEAN VALUE THEOREM AND DARBOUX'S PROPERTY FOR THE Derivative of an additive set function with respect to a measure on eⁿ

by R. J. Easton, and S. G. Wayment

1) INTRODUCTION: During a recent investigation of existence and equality almos everywhere of the cross partial derivatives f_{xy} and f_{yx} , a somewhat different derivative Df for a function f(x,y) was used (4). This derivative Df is also defined for a function f of n-variables. The purpose of this paper is to establish a mean val ue theorem and a Darboux property for the function Df, and to generalize these results to the derivative T'* as defined on pp. 268-271 of (8). A similar result was obtained by L. Misik in (7), but the technique is much more cumbersome. The method of proof is the same as that given in (5). In the case n=2 the technique of proof is used to establish a theorem concerning the equality of the three derivatives. For simplicity the proofs and definitions will be given for n=2.

II) THE DERIVATIVE Df: Let $R = [a,b;c,d] = \{(x,y) | x \in [a,b], y \in [c,d] \}$. If f(x,y) is a function whose domain contains R, then the f-area of R is denoted by F(R) = f(b,d) - f(a,d) - f(b,c) + f(a,c). The ordinary area of R will be denoted by A(R). A rectangle R = [a,b;c,d] is said to be of order M, if M > 1 and 1/M < (d-c)/(b-a) < M. One then defines the upper and lower derivatives of order M at a point (x,y) to be $\overline{\lim}$ and \lim respectively of ratios of f-areas to ordinary areas of rectangles or order M which contain (x,y) and whose areas converge to zero. Then f is said to have a derivative of order M, $D_M f(x,y) = D_M f(P)$, at P = (x,y) if the upper and lower derivatives of order M are equal. The function f is said to be two nondecreasing if the f-area of each sub-rectangle of R is non-negative. It follows {2}, that if f is of bounded variation in the sense of Hardy, then except for a set of measure zero, $D_{N}f = D_{M}f$ for each $N \ge M \ge 1$. The common value is denoted by Df. It also follows {4}, that f_{xy} and f_{yx} each exist except possibly on a set of measure zero.

III) THE MEAN VALUE THEOREM:

THEOREM 1: If $D_M f$ exists at each point P of a closed rectangle R of order M and f is continuous at each point of R, then there exists a point Q ϵ R such that $D_M f(Q) = F(R)/A(R)$.

Proof: Suppose $R_1 = R = [a,b;c,d]$ and divide R_1 into four rectangles using the lines x = a+h/2, y = b+k/2 where h = b-a and k = d-c. Denote the rectangles, beginning in the lower left hand corner and proceeding counterclockwise, by R_{11} , R_{12} , R_{13} , R_{14} , and observe that each of the four rectangles is similar to R_1 . It follows that

$$\sum_{i=1}^{4} F(R_{1i}) = F(R_1) \text{ and } \sum_{i=1}^{4} A(R_{1i}) = A(R_1)$$

and hence there must exist a j and a k such that

$$F(R_{1i}) \ge (1/4)F(R_1)$$
 and $F(R_{1k}) \le (1/4)F(R_1)$.

We now proceed to find a rectangle R_2 of order M with sides parallel to the sides of R_1 such that $R_2 \subseteq R_1$, $F(R_2) = (1/4)F(R_1)$, and $A(R_2) = (1/4)A(R_1)$. If $F(R_{1i}) = (1/4)F(R_1)$ for some i, then choose $R_2 = R_{1i}$. Suppose equality does not hold for any i and consider the case j=3 and k=1. The other cases would follow in a similar manner. Let $\alpha = h/k$ and define the auxiliary function

 $g(t)=f(a+t+h/2,b+\alpha t+k/2)-f(a+t,b+\alpha t+k/2)-f(a+t+h/2,b+\alpha t)+f(a+t,b+\alpha t).$

Then $g(0) = F(R_{11})$, $g(h/2) = F(R_{13})$, and g is a continuous function of t for $0 \le t \le h/2$. The ordinary intermediate value theorem for a function of one variable guarantees the existence of a $t_0 \varepsilon(0, h/2)$ such that $g(t_0) = (1/4)F(R_1)$. This value t_0 defines R_2 and we note that $F(R_2)/A(R_2) = F(R_1)/A(R_1)$. In the sequel we shall refer to the above selection process for determining R_2 as the sliding technique. We proceed inductively to define a nested sequence of closed rectan gles $\{R_i\}$, each of order M with sides parallel to R_1 , such that

i)
$$F(R_{i+1}) = (1/4)F(R_i)$$
 and

ii)
$$A(R_{i+1}) = (1/4)A(R_i)$$
.

By the nested interval theorem there exists exactly one point $Q \in \bigcap R_i$, and

$$D_{M}f(Q) = \lim F(R_{1})/A(R_{1}) = F(R_{1})/A(R_{1}).$$

We shall say that the set function F has property I provided the auxiliary function g(t) has the intermediate value property along the lines x = constant, y = constant, and y = tax. We have the somewhat stronger result.

THEOREM 2. If $D_M f$ exists at each point P of a closed rectangle R of order M and the set function F has property I, then there exists a point Q ϵ R such that $D_M f(Q) = F(R)/A(R)$.

The following example shows that theorem 2 is a stronger result.
EXAMPLE 1: Let f(x,y) be defined as follows on the unit square:

f(x,y) = 1 if y is rational and

f(x,y) = 0 if y is irrational

Then $D_M f$ exists and is zero at each point of the unit square and the set function F has property I.

We now proceed to remove the condition that $D_M f$ exist along the bound ary of R.

THEOREM 3. If R is a rectangle of order M, $D_M f$ exists at each point P ϵ int(R), and F has property I, then there exists a point $Q\epsilon$ int(R) such that $D_M f(Q) = F(R)/A(R)$.

Proof: We will use the same notation as in theorem 1 and modify the selection process to obtain, for some k, an $R_k \subset int(R)$. We consider the following two cases:

- 1) If $R_2 \neq R_{1i}$ for any i, then R_2 has at most one edge contained in bdry(R).
- If R₂ = R_{1i} for some i, say i=1, then the following argument allows us to choose R₃ with at most one edge contained in bdry(R).

Divide R_2 into R_{21} , R_{22} , R_{23} , R_{24} , and if $F(R_{21}) = F(R_{22}) = F(R_{23}) = F(R_{24}) = (1/4)F(R_2)$, choose $R_3 = R_{23} \subset int(R)$. If $F(R_{21}) \neq (1/4)F(R_2)$ for some i, then the sliding technique gives an R_3 with at most one edge contained in bdry(R). For case (1), suppose the bottom edge of R_2 is contained in bdry(R_1) and divide R_2 as before. If $F(R_{21}) = (1/4)F(R_2)$ for i=3, choose $R_3 = R_{23}$ and if $F(R_{23}) \neq (1/4)F(R_2)$, then the sliding technique will again give an $R_3 \subset int(R)$.

IV) A DARBOUX PROPERTY:

THEOREM 4: Let \emptyset be a connected open set in E^2 . Suppose $D_M f$ exists at each point of \emptyset and that F has property I. Let P, $Q \in \emptyset$ and suppose $D_M f(P) = \alpha$, $D_M f(Q) = \beta$, $\alpha < \beta$, and $\lambda \in (\alpha, \beta)$. If \overline{PQ} is an arc which is contained in \emptyset with endpoints P and Q, then for each $\varepsilon > 0$ there exists a point $S \in \emptyset$ such that the distance $d(S, \overline{PQ})$ from S to \overline{PQ} is less than ε and $D_M f(S) = \lambda$.

Proof: Construct a polygonal arc PQ from P to Q consisting of horizontal and vertical straight line segments such that each point of PQ is within min($\xi/2, \varepsilon/2$) of \overline{PQ} , where $\xi = d(\overline{PQ}, bdry 0)$. Let u =

= min $(\epsilon/2, (\lambda-\alpha)/2, (\beta-\lambda)/2)$. There exist rectangles R₁ and R₂ of order M centered at P and Q respectively with edges parallel to the coordinate axes and having the same base and height, such that

 $|F(R_1)/A(R_1)-\alpha| < u$, $|F(R_2)/A(R_2)-\beta| < u$, and diam $(R_1) = diam(R_2) \le s \le min(\varepsilon/2, \varepsilon/2)$.

Since F has property I the sliding technique allows us to obtain a rectangle R or order M, with sides parallel to the coordinate axes, centered at a point of \hat{PQ} such that diam(R) = diam(R₁) and F(R)/A(R)= = λ . Theorem 1 implies the existence of a point S ϵ R such that $D_{M}f(S) = \lambda$ and $d(S, \overline{PQ}) < \epsilon$.

REMARK. Let u be Lebesgue measure on E^n and let T be any absolutely continuous measure with respect to u. Further suppose that the derivative T' *(x), as defined in {8}, exists at every point in an in terval $R_n \subset E^n$. The above technique may be used to establish a mean value theorem and a Darboux property for this derivative. These results also hold if T is an additive set function defined on at least the closed intervals in E^n and has the intermediate value property along straight lines in the appropriate directions, u is a translational invariant measure which is finite on regular rectangles, and T' *(x) exists at every point in Int (R). A further generalization is given in Section VI of this paper.

V) A THEOREM ON THE EQUALITY OF THE DERIVATIVES f_{xv} , f_{vx} , and Df.

It is well known that if $f_{xy}(x,y)$ exists at each point of an open set \emptyset and R is a closed rectangle with $R \subset \emptyset$, then there exists a point P ε int (R) such that $f_{xy}(P) = F(R)/A(R)$. Example 1 shows that there are functions for which $f_{xy}(x,y)$ and Df exist on a rectangle and f_{yx} fails to exist at any point. Also, the example can be modified by defining f(x,y) = 2 whenever x and y are rational, f(x,y) = 1 if exactly one of x or y is rational, and zero otherwise to give a function such that Df exists on a rectangle while both f_{xy} and f_{yx} fail to exist at any point.

THEOREM 5. If f_{xy} and $D_M f$ exist on an open set 0 and (a) the function f_{xy} is continuous or (b) the related set function F has the intermediate value property along straight lines in the appropriate directions and $D_M f$ is continuous, then $f_{xy}(P) = D_M f(P)$ for each $P \in 0$.

Proof: Suppose f_{xy} is continuous. Let $\{R_i\}$ be a sequence of nested

rectangles of order M contained in 0 and closing down on P. Then $D_M f(P) = \lim_{i \to \infty} F(R_i)/A(R_i)$. For each i there exists a point $P_i \in R_i$ such that $f_{xy}(P_i) = F(R_i)/A(R_i)$ and the continuity of f_{xy} gives the desired result.

Suppose that $D_M f$ is continuous and that the set function F has the intermediate value property and hence theorem 2 applies. Let $P = (x_o, y_o) \in \mathcal{O}$. If $\varepsilon > 0$ then there exists t_1 such that $0 < t_1 < \varepsilon$ and $|[f_x(x_o, y_o+t_1)-f_x(x_o, y_o-t_1)]/2t_1-f_{xy}(x_o, y_o)| < \varepsilon/3$. There exists $t_2 > 0$ such that $t_2 = nt_1$ for some integer n and so that $[[f(x_o+t_2, y_o+t_1)-f(x_o-t_2, y_o+t_1)/[2t_2]-f_x(x_o, y_o+t_1)]/[2t_1] < \varepsilon/3$ and $[[f(x_o+t_2, y_o-t_1)-f(x_o-t_2, y_o-t_1)/[2t_2]-f_x(x_o, y_o-t_1)]/[2t_1] < \varepsilon/3$. We can now divide rectangle $R = [x_o-t_2, x_o+t_2; y_o-t_1, y_o+t_1]$ into n squares and conclude from the sliding technique that there exists a square R' such that $R' \subset R$ such that $D_M f(P') = F(R)/A(R)$ and hence $|D_M f(P')-f_x(P_0)| < \varepsilon$ and $d(P', P) < \varepsilon$. Continuity implies the desired result.

VI) A FURTHER GENERALIZATION:

We shall say that the additive set function F has property C provided the auxiliary function g(t) as defined in theorem 1 is continuous along the lines x=constant, y=constant, and y = $\pm \alpha x$. Note that property C implies property I.

THEOREM 6. Suppose that S and T are additive set functions defined on rectangles, u is a translational invariant measure, and S' *(p) and T' *(p) exist at each point $p \in Int$ (R) and T' *(p) $\neq 0$ for any p. Then there exists a point $q \in Int$ (R₂) so that

$$\frac{S(R_{0})}{T(R_{0})} = \frac{S' * (q)}{T' * (q)}$$

Proof: Let $U(R) = S(R_0)T(R)-T(R_0)S(R)$. Then $U(R_0) = 0$ and U has property I. Hence there exists a point $q \in Int(R_0)$ so that U' *(q)= = 0 = $S(R_0)T'$ *(q)- $T(R_0)S'$ *(q). This holds without the condition that T' *(p) $\neq 0$ for $p \in R_0$. The result now follows.

The method of proof in the preceding theorem allows one to remove the condition of translational invariant u. Suppose S and T are additive

set functions defined on rectangles and having property C. Define dS(p)/dT to be the limit of the ratio of the S area to the T area of regular rectangles as the diameters of the rectangles tend to zero.

THEOREM 7. If dS/dT exists at each point p in Int (R_0) , $T(R) \neq 0$ for $R \subset R_0$, and S and T have property C, then there is a point $q \in Int (R_0)$ so that $dS(q)/dT = S(R_0)/T(R_0)$.

Proof: Define U as in theorem 6. Then use the procedures of theorems 1 and 3 to define a nested sequence $\{R_i\}$ of rectangles closing down on q ε Int R_o and such that $U(R_i) = 0$. Then $S(R_o)/T(R_o) =$ = $S(R_i)/T(R_i)$ and the result follows.

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REVISTA DE LA UNION MATEMATICA ARGENTINA Volumen 24, Número 3, 1969

NOTES ON COMARGINAL PROBABILITY MEASURES

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1. INTRODUCTION.

By $S = (\Omega, A, B, C, P_1, P_2)$ we shall denote a system consisting of two σ -algebras B, C of subsets of Ω , the generated σ -algebra $A=\tau(B,C)$, and probabilities P_1 , P_2 defined on B,C respectively, which are compatible in the sense that $P_1 = P_2 = P$ on the intersection σ -algebra $\mathcal{D} = B \cap C$.

A probability Q on A such that its restrictions to B, C are P_1, P_2 , respectively, is said to be *comarginal* with P_1 , P_2 , or briefly comarginal.

A probability Q on A, not neccesarily comarginal, is called *commutative* if on A-measurable Q-integrable functions:

(1) $E_Q^B \cdot E_Q^C = E_Q^C \cdot E_Q^B = E_Q^D$ [Q],

where E_Q^G denotes the conditional expectation operator with respect to the σ -algebra $G \subset A$ and to the measure Q.

The main problem we consider here is to search under'what conditions a system S admits a comarginal and commutative measure Q. For such a measure we can assert its uniqueness. Owing to this fact and to the following example we shall call Q the (generalized) product measure on S.

Given the probability spaces (Ω_1, B', P_1') , (Ω_2, C', P_2') the system S formed by $\Omega = \Omega_1 \times \Omega_2$, $B = \pi_1^{-1}(B')$, $C = \pi_2^{-2}(C')$, $A = \tau(B,C)$, $P_i = P_i' \pi_i$, $(\pi_i = \text{projection on } \Omega_i)$, i=1,2, admits the comarginal and commutative measure $Q = P_1' \times P_2'$. The relation of commutation (1) is here Fubini's theorem.

In this example $\mathcal{D} = \{\phi, \Omega\}$; the opposite extreme case, when $\mathcal{D} = B$ (or $\mathcal{D} = C$), gives us also an example very trivial for a product measure: $Q = P_2$ (or $Q = P_1$).

Intermediate cases can be given as follows:

Let $\Omega = \Omega_x \times \Omega_y \times \Omega_z$ and probabilities dx, dy, dz be given on σ -algebras B_x , B_y , B_z of Ω_x , Ω_y , Ω_z respectively. We define $B = \{\phi, \Omega_x\} \otimes B_y \otimes B_z$, $C = B_x \otimes \{\phi, \Omega_y\} \otimes B_z$ and the probabilities $dP_1 = dy dz$, $dP_2 = dx dz$. Here \mathcal{D} is isomorphic to B_z and the system $(\Omega, A, B, C, P_1 P_2)$ has the product measure dQ = dx dy dz.

The preceding situation always appears in a Markov process. Assume

 $\{x_i\}_{i=1,2..}$ is a Markov chain and $B = \tau(x_1, \ldots, x_n)$, $C = \tau(x_n, x_{n+1}..)$ and $\mathcal{D} = \tau(x_n)$. The product measure Q turns out to be the probability associated with the process, and the relation of commutation (1) can be rewritten as expressing the conditional independence of past (B) and future (C), given the present (p). (cf. [1], \$14, or [4], part A, ch. II).

For further examples we refer to $\{1\}$, §3.

2. COMARGINAL MEASURES AND BILINEAR FORMS.

We shall say that T: $B \times C \longrightarrow [0, +\infty)$ is a (positive) bilinear form on the system $S = (\Omega, A, B, C, P_1, P_2)$ if T(B,C) is additive in each variable separately; and T will be said to be *compatible* (with P_1, P_2) if T(B, Ω) = P₁(B), T(Ω, C) = P₂(C), for any B \in B, C \in C.

To each finitely additive measure Q on the algebra $A_0 = B \vee C \operatorname{comar}$ ginal with P_1 , P_2 we associate the compatible bilinear form T(B,C) = Q(BC), which verifies: T(B,C) = 0 if $BC = \emptyset$. Conversely:

THEOREM 1. Every compatible bilinear form T on C such that T(B,C)=0 whenever BC = ϕ , defines on A₀ a unique finitely additive comarginal probability given by Q(B,C) = T(B,C).

Proof: We define Q(BC) = T(B,C). If BC = B'C', then $BC = (BB')(CC') = B_1C_1$. To prove that T(B,C) = T(B',C') it is enough to prove that $T(B,C) = T(B_1,C_1)$. It follows from $T(B,C) = T(B_1,C) + T(B-B_1,C) = T(B_1,C)$, since $(B-B_1)C = \emptyset$, and from $T(B_1,C) = T(B_1,C_1)$. Hence Q is well defined on sets of the form BC, $B \in B$, $C \in C$.

Let BC = $\sum_{\alpha} B_{\alpha}C_{\alpha}$, where α runs on a finite family of indices. In or der to prove that Q can be (uniquely) extended to A_{0} as a finitely additive measure it is enough to prove that Q(BC) = $\sum_{\alpha} Q(B_{\alpha}C_{\alpha})$.

Let $\{B_i\}$ be the partition of B defined by the B_{α} 's and $\{C_j\}$ that defined by the C_{α} 's on C. We can assume from the beginning that $B_{\alpha} \subset B$, $C_{\alpha} \subset C$. Since $B_{\alpha} = \sum_{m} B_{\alpha,m}$, $C_{\alpha} = \sum_{n} C_{\alpha,n}$, denoting by $B_{\alpha,n}$ ($C_{\alpha,n}$) the sets of the mentioned partition of B (C) included in B_{α} (C_{α}), we have:

$$\sum_{i,j} B_i C_j = BC = \sum_{\alpha} B_{\alpha} C_{\alpha} = \sum_{\alpha} (\sum_{m} B_{\alpha,m}) (\sum_{n} C_{\alpha,n}) = \sum_{\alpha} (\sum_{m,n} B_{\alpha,m} C_{\alpha,n})$$

This means that in the first and last sums appear the same non-void terms. Therefore, from the bilinearity of T we get

 $Q(BC) = T(B,C) = \sum_{i,j} T(B_i,C_j) = \sum_{\alpha} (\sum_{m,n} T(B_{\alpha,m},C_{\alpha,n})) = \sum_{\alpha} T(\sum_{m} B_{\alpha,m},\sum_{n} C_{\alpha,n}) =$

REMARKS. 1) It is not true, in general, that Q is σ -additive. In fact, let Ω be the triangle on the plane defined by x > 0, y > 0, x + y < 1. B (C) the Borel sets of Ω depending only of x(y). For B ε B (C ε C) we define P₁(B) = m(B'), (P₂(C) = m(C')), where B'(C) denotes the projection of B (C) on the x (y) axis and m the Lebesgue measure on (0,1).

The bilinear form $T(B,C) = \int_{0}^{1} 1_{B'}(t) 1_{C'}(1-t)dt$ is compatible with P_1 , P_2 but, as it is easy to see, Ω can be put as a countable sum of rectangles B.C for which T(B,C) = 0.

2) Let us observe that if T is bilinear and comarginal then T is also σ -bilinear; i.e. σ -additive in each variable. In fact,

$$T\left(\sum_{j=0}^{\infty} B_{j}, C\right) = T\left(\sum_{i=1}^{n} B_{j}, C\right) + T\left(\sum_{n+1}^{\infty} B_{j}, C\right), \text{ then}$$
$$\left|T\left(\sum_{j=1}^{\infty} B_{j}, C\right) - \sum_{i=1}^{n} T\left(B_{j}, C\right)\right| \leq \sum_{n+1}^{\infty} T\left(B_{j}, \Omega\right) = \sum_{n+1}^{\infty} P_{1}\left(B_{j}\right) \longrightarrow 0 \text{ if } n \longrightarrow \infty.$$

Then the proof of the theorem remains true if we assume that α runs on a countable family of indices such that the B_{α} 's and the C_{α} 's de fine, respectively, countable partitions of the spaces α .

For example, we can assert the $\sigma\text{-additivity}$ of Q if A is defined by a countable partition of $\Omega.$

3) In theorem 1 we can assume $T(B_o, C_o) = 0$ if $B_oC_o = \emptyset$ for $B_o\varepsilon B_o C_o = B_o C_o = C$, where B_o, C_o are collections of sets with the *approximation property*: $P_1(B) = \sup_{b_o \in B_o} P_1(B_o)$, $P_2(C) = \sup_{c_o \in C} P_2(C_o)$, for $B \in B, C \in C$. In fact, if $B_o = B$, $C_o = C$, and $B.C = \emptyset$: $T(B,C) = T(B-B_o+B_o, C-C_o+C_o) = T(B-B_o, C) + T(B_o, C-C_o) + T(B_o, C_o) \le P_1(B-B_o) + P_2(C-C_o)$. The last member can be done arbitrarily small, hence T(B,C) = 0.

4) The σ -additivity of Q follows under the following hypothesis

1) $K_B \subset B$, $K_C \subset C$ are *semi-compact classes* (i.e. every countable family of K_B (K_C) with an empty intersection has a finite subfamily which also intersects in the empty set) verifying the approximation property (as defined above).

2) $K_B \cdot K_C = \{K.L; K \in K_B, L \in K_C\}$ is a semicompact class of sets. In fact, the class *L* of finite unions of sets of $K_B \cdot K_C$ enjoys the property of approximation in $B \lor C$, since, as it was shown above, for $K \in B$, $L \in C$ we have $Q(BC - KL) \leq P_1(B - K) + P_2(C - L)$, then Q(BC - KL) can be done arbitrarily small. *L* being compact

0.E.D.

and with the approximation property the σ -additivity of Q follows from a theorem of Alexandrov (cf {4}, pp. 47).

3. ∇-COMMUTATIVE SYSTEMS.

For any probability space (Ω, A, P) we define the measurable hull of $X \subset \Omega$ as a set $A \in A$ containing X except by a set of P-outer measure zero and with minimal P-measure. Of course, the measurable hull is defined except on a null set of A, and it provides a well defined <u>e</u> lement of the Boolean measure algebra: A/[P]. If $B \subset A$ is a σ -sub-algebra of A the measurable hull of X ϵ A with respect to B coincides with $\{E_{p}^{B} | 1_{v} > 0\}$, [P].

For a system $S = (\Omega, A, B, C, P_1, P_2)$ we shall designate $\nabla^1 X$, $\nabla^2 X$, ∇X the measurable hulls of $X \subset \Omega$ with respect to B, C, D and to the measures P_1 , P_2 , P respectively.

If Q is a comarginal measure on S and E, F, G denote E_Q^B , E_Q^C , E_Q^p respectively, we can see that the condition EF = FE = G[Q] (on boun ded A- measurable functions) is equivalent to Ef = Gf $[P_1]$ (on C-measurable functions) and also to Ff = Gf $[P_2]$ (on B-measurable functions) (cf {1}, \$ 1).

For a comarginal measure Q on S, the condition EF = FE = G implies $\nabla^1 \nabla^2 X = \nabla^2 \nabla^1 X = \nabla X$ [Q] for X ε A. This condition is equivalent to $\nabla^2 B = \nabla B$ [P₂], for any B ε B, and also to $\nabla^1 C = \nabla C$ [P₁], for any C ε C. (cf. {1}, \$5).

We note that any of these last conditions can be introduced even if we do not assume that a comarginal measure Q is known. Then we adopt the following definition:

We shall say that the system $S \nabla$ -commutes if $\nabla^2 B = \nabla B [P_2]$, $\forall B \in B$. From the above considerations it follows:

In order that there exists a comarginal and commutative measure on $S = (\Omega, A, B, C, P_1, P_2)$ it is necessary that S be a ∇ -commutative system.

An independent proof will be given in next theorem 2.

A ∇ -commutative system is said to be *simple* if $\mathcal{D}/[P]$ is the Boolean algebra $\{0,1\}$.

4. THE FINITELY ADDITIVE MEASURE ASSOCIATED TO A ∇ -commutative system.

Given the system $S = (\alpha, A, B, C, P_1, P_2)$ we observe that the conditional expectation operator G can be calculated on B (C)-measurable functions, with respect to P_1 (P_2) even if there is no comarginal measure. Hence we can define the compatible bilinear form:

(2)
$$T(B,C) = \int GI_B \cdot GI_C dP$$
.

THEOREM 2. In order that Q(BC) = T(B,C) defines a finitely additive comarginal measure on S it is neccessary and sufficient that S be a ∇ -commutative system.

In this case, if Q is a probability on A (i.e. if Q is σ -additive), Q is the unique commutative comarginal measure on S. (uniqueness of the product measure).

Proof: By theorem 1, to prove that ∇ -commutativity implies that Q is a finitely additive measure, we have to show that B.C = \emptyset implies T(B,C) = 0. From B.C = \emptyset we get $\nabla^2 B.C = \nabla B.C = \emptyset[P_2]$; i.e. $\{G1_B > 0\}.C =$ $= \emptyset[P_2]$. Then $T(B,C) = \int_C G1_B.dP_2 = 0$. Conversely, if Q is a finite ly additive and comarginal measure, S ∇ -commutes. In fact, $\nabla^2 B$ $C \nabla B[P_2]$ and on the other hand from $P_2^*(B - \nabla^2 B) = 0$ we have $B - \nabla^2 B$ $C' \in C$ with $P_2(C') = 0$. Hence $Q(B - \nabla^2 B) \leq Q(C') = P_2(C') = 0$. Then, $Q(B - \nabla^2 B) = \int_{C} G1_B.dP_2 = 0$. This means $\{G1_B > 0\}.(\nabla^2 B = \emptyset[P_2])$ which implies $\nabla B \subset \nabla^2 B[P_2]$. If Q, as defined above, is a measure, and E, F the conditional expectations with respect to B.C :

(3)
$$Q(B.C) = \int G_B \cdot G_C dP = \int_C G_B dQ = \int_C F_B dQ$$

Then $G_{B}^{1} = F_{B}^{1}$, $V = E_{\epsilon}^{1}$ B. This implies the commutation of Q. Another comarginal commutative measure Q' must verify (3), but since

$$\int_{C} G_{B}^{1} dQ' = \int_{C} F_{B}^{1} dQ' = Q'(B.C) = Q(B.C) ,$$

Q and Q' coincide on $B \lor C$, and therefore on A. QED.

REMARK: The condition $\nabla^2 B = \nabla B [P_2]$, $B \in B$, defining a ∇ -commutative system implies the symmetric one $\nabla^1 C = \nabla C [P_1]$, $C \in C$.

In fact, the first one implies that $\int G_{B}^{1} G_{C}^{1} dP$ defines a finitely ad ditive measure, and from this we derive $\nabla^{1}C = \nabla C [P_{1}]$ as it was done with $\nabla^{2}B = \nabla B [P_{2}]$ in the proof of theorem 2. Now, we obtain easily

 $\nabla^2 \nabla^1 X = \nabla^1 \nabla^2 X = \nabla X [P]$ for each $X \in B \vee C$.

THEOREM 3. i) If the ∇ -commutative system S is simple (i.e. $D/[P] = \{0,1\}$) then $Q(B.C) = P_1(B) \cdot P_2(C)$. In particular for a system S obtained from a cartesian product (as described in the introduction), $Q = P'_1 \times P'_2$.

ii) If, conversely, $Q(B.C) = P_1(B) \cdot P_2(C)$ defines a finitely additive measure on $B \lor C$, S is a simple ∇ -commutative system.

Proof: i) It follows from theorem 2 and the fact that $G1_B = P_1(B) \cdot I_{\Omega}$, $G1_C = P_2(C) \cdot I_{\Omega}$.

ii) If $D \in \mathcal{D}$ we have $Q(D) = P(D) = Q(D.D) = P(D)^2$, then P(D) = 0 or 1. $G1_B$, $G1_C$ are computed like in i), then $\int G1_B.G1_C = P_1(B).P_2(C)$. From theorem 2) it follows that S is a ∇ -commutative system.

We shall say that a system S is *complete* with respect to a comarginal probability P defined on A if every P-null set of A belongs to P.

THEOREM 4. If the system S is complete with respect to a comarginal probability P and R is a commutative probability on S equivalent to P then the product measure Q exists and it is equivalent to P.

Proof: Assume $f = \frac{dR}{dP}$. By hypothesis $F_R = G_R$ on B-measurable functions and $E_R = G_R$ on C-measurable functions. We have (c.f. { 1 } §2):

$$E_{R}(h) = E(f.h) / E(f)$$
,
 $F_{R}(g) = F(f.g) / F(f)$, $G_{R}(m) = G(f.m) / G(f)$,

where E, F, G denote here the conditional expectation operators with respect to the measure P and B, C, D respectively.

In {1}, th. 2, \$10, it is proved that the probability measures equivalent to R that also commute are characterized as those whose Radon-Nikodym derivatives with respect to R are of the form: g.h, where g (h) is a B(C)-measurable function (both positive and finite [R]). Let us consider the functions:

$$g = \frac{+\sqrt{Gf}}{Ef}$$
 , $h = \frac{+\sqrt{Gf}}{Ff}$

Then, g.E(h.f) = Gf.E(f/Ff) /Ef = Gf.E_R(1/Ff) = Gf.G_R(1/Ff) = G(f/Ff) = = GF(f/Ff) = 1. Analogously h.F(g.f) = 1. From $\int_{\mathbb{R}} ghf dP = \int_{\mathbb{R}} gE(hf) dP =$ $= \int_{B} 1 \, dP = P(B) \text{ and } \int_{C} ghf \, dP = P(C)$ We see that the probability Q defined by

 $Q(A) = \int_{A} gh.f dP = \int_{A} gh dR$ is comarginal with P and since $\frac{dQ}{dR} = g.h$, it commutes and obviously $Q \sim R$. QED.

REMARK: If the system S is complete with respect to P any other comarginal measure R is absolutely continuous with respect to P, since P(A) = 0 implies A εP and then R(A) = P(A) = 0. Then, if a product measure Q exists on S, we have Q \ll P. In spite of the fact that for B ε B, C ε C, Q(B.C) = 0 implies P(B.C) = 0 (since $\int G1_BG1_CdP = 0$ implies $P_2(\nabla B.C) = 0$ and then P(B.C) = 0) we have not Q $\sim P$, in general. Let us see the following example:

Let Ω be the product of X and Y, X = Y = (0,1), A = the Borel sets of X × Y, B (C) the Borel sets independent of y (x), P the probability on A equal to $\frac{m \times m}{2}$ (m the Lebesgue measure on (0,1)) plus a measure of total mass 1/2 concentrated on the diagonal and uniformely distributed there.

The product measure $Q = m \times m$ is not equivalent to P.

If $D \in \mathcal{D}$ defines an atom of the σ -algebra $\mathcal{D}/[P]$ it can be seen that $S_{\rm D} = (D, A \land D, B \land D, C \land D, P_1/P_1(D), P_2/P_2(D))$ is a simple ∇ -commutative system, whenever S is a ∇ -commutative system. Moreover, if Q is the (product) finitely additive measure on S defined above , $\frac{Q(B.C.D)}{Q(D)}$ is the product measure on $S_{\rm D}$ as it is easily seen. Since $S_{\rm D}$ is simple: $\frac{Q(B.C.D)}{Q(D)} = \frac{P_1(B.D)}{P(D)} \cdot \frac{P_2(C.D)}{P(D)}$, and then

(4)
$$Q(B.C.D) = P_1(B.D).P_2(C.D)/P(D)$$

THEOREM 5. If in a ∇ -commutative system S, A is defined by a countable partition of Ω , there exists the product measure Q. If $\{D_i\}$ is the partition defining D, then Q is defined by:

(5)
$$Q(B.C) = \sum_{i} \frac{P_1(BD_i) \cdot P_2(CD_i)}{P(D_i)}$$

Proof: The σ -additivity of Q follows from the second remark in §2. We have Q(B.C) = $\sum_i Q(B.C.D_i)$, and applying (4) we obtain the equality (5).

5. THE BOOLEAN MEASURE STRUCTURE OF A V-COMMUTATIVE SYSTEM

From a given probability space (Ω, A, P) and a σ -subalgebra B of A, we get the (measure) Boolean algebra A = A/[P] quotient of A mod. P-null sets and the subalgebra B = B/[P]. To the operation of take ing the B-measurable hull in A corresponds in A a so called monadic operator (c.f. {2}).

Let A be a Boolean algebra with a subalgebra B such that for each a ε A there exists an element b ε B which is the least element of B verifying b \geq a ; we set b = ∇ a, and we call ∇ the monadic operator in A related to B.

A monadic operator verifies the following properties: $\nabla 0 = 0$, $\nabla(a \lor b) = \nabla a \lor \nabla b$, $\nabla \nabla a = \nabla a$, $\nabla(a \land \nabla b) = \nabla a \land \nabla b$.

The algebraic system (A, B, ∇) is called a monadic algebra. Let us consider a ∇ -commutative system $S = (\Omega, A, B, C, P_1, P_2)$ with a product measure Q. By passing to the quotient mod. [Q] we get the Boolean (measure) algebras A, B, C, D, which are the images of A, B, C, D resp.; B, C, D are subalgebras of A and D = B \cap C. If we denote by $A_0 = B \vee C$ then $A_0 = B \vee C \simeq A_0/[Q]$, where $B \vee C$ is the Boolean algebra generated by B and C.

If $\nabla^1, \nabla^2, \nabla$ designate the monadic operators in A corresponding to the measurable hull operations in A denoted before with the same symbols we have $\nabla^1 \nabla^2 a = \nabla^2 \nabla^1 a = \nabla a$ for any $a \in A$. The same hold if we restrict our-selves to elements $a \in A_o = B \vee C$. So we have an instance of what is called a biadic algebra (c.f. {2}).

The algebraic system $(A_o, B, C, \nabla^1, \nabla^2)$ is called a *biadic algebra* if (A_o, B, ∇^1) , (A_o, C, ∇^2) are monadic algebras, $A_o = B \lor C$ and $\nabla^1 \nabla^2 = \nabla^2 \nabla^1$. It is easy to see that in this case $\nabla = \nabla^1 \nabla^2$ defines the monadic operator related to $D = B \cap C$.

We can say that the underlying Boolean structure of a ∇ -commutative system is a biadic algebra. We have seen this when a product measure Q is given in the system and it is easy to see that the same is true even if Q does not admit a σ -additive extension from A_c to A.

In particular, if S is a simple ∇ -commutative system, we obtain a simple biadic algebra $(A_0, B, C, \nabla^1, \nabla^2)$ i.e. it verifies $D = B \wedge C =$ = {0,1}. For such simple algebras we have $A_0 = B \oplus C$, direct sum of B, C which means that $A_0 = B \vee C$ and, for $b \in B$, $c \in C$, $b \wedge c = 0$ implies b = 0 or c = 0 (in fact, if $b \wedge c = 0$ then $0 = \nabla(b \wedge c) =$ = $\nabla^1 \nabla^2 (b \wedge c) = \nabla^1 (\nabla^2 b \wedge c) = \nabla^1 (\nabla^2 \nabla^1 b \wedge c) = \nabla^1 (\nabla^2 \nabla^1 b \wedge c) = \nabla b \wedge \nabla c$, therefore $\nabla b = 0$ or $\nabla c = 0$, so b = 0 or c = 0 (The converse also holds: if $A_0 = B \oplus C$, $(A_0, B, C, \nabla^1, \nabla^2)$ is a simple biadic algebra).

In simple ∇ -commutative systems, for example the systems obtained from cartesian products, what really matters from the point of view of the theory of measure Boolean algebras are the algebras B, C and A_o . Explicitly, if A'_o , B', C' are obtained from another simple ∇ -commutative system S', then if we have Boolean isomorphisms B = B' and C = C' we get $A_o = A'_o$. This is due to the fact that A_o , A'_o are direct sums. In fact, the direct sum $A_o = B \oplus C$ has the property of extension of homomorphisms: if $B \longrightarrow \bar{A}$ and $C \longrightarrow \bar{A}$ are Boolean homomorphisms with range a Boolean algebra \bar{A} , there exists one and only one extension of them to a homomorphism: $A_o = B \oplus C \longrightarrow \bar{A}$. (c.f. {5}).

On the other hand, if (B,P_1) (C,P_2) are given Boolean measure algebras we can construct at least a simple ∇ -commutative system S for which the associated biadic algebra is precisely $(B \oplus C, B, C, \nabla^1, \nabla^2)$. In fact, it is well known that the Stone space of $B \oplus C$ is the cartesian product of the Stone spaces of B and C, $S(B \oplus C) = S(B) \times S(C)$ (Precisely the algebra of clopens of $S(B) \times S(C)$ is used to define $B \oplus C$). We set on the clopens of S(B) and S(C) the measures P_1 and P_2 in the obvious way and we extend P_1 , P_2 to the σ -algebras generated by clopens associated to elements of B, C respectively. The product of the probability spaces so obtained gives us the required system S.

For general ∇ -commutative systems we can prove analogous results.

To a ∇ -commutative system $S = (\Omega, A, B, C, P_1, P_2)$ we have associated a biadic algebra (A_0, B, C) . Moreover B,C are measure Boolean algebras with the probabilities P_1 , P_2 defined on B, C, respectively, coinciding in $D = B \cap C$. Let us call $M(S) = (A_0, B, C, P_1, P_2)$ this Boolean measure structure associated with S. We shall say that the ∇ -commutative systems S, S' have the same Boolean measure structure, $M(S) \equiv M(S')$, if under a unique Boolean isomorphism $A_0 \cong A'_0$, $B \cong B'$, $C \cong C'$ and $D \cong D'$; and the probabilities P_1 , P_1' (i=1,2) correspond under the isomorphism.

THEOREM 6. 1) Given two probability Boolean algebras $(B,P_1), (C,P_2)$ and sub-s-algebras $D \subset B$, $D' \subset C$ such that $D \cong D'$ under a fixed isomorphism preserving the measures $P_1 | _D$, $P_2 | _D'$, there exists a ∇ -com mutative system S such that $M(S) = (\bar{A}_0, \bar{B}, \bar{C}, \bar{P}_1, \bar{P}_2)$, where $\bar{B} \cong B, \bar{C} \cong C$ are measure preserving isomorphisms (with respect to P_1 , $\bar{P}_1 = 1, 2$) which restricted to $\bar{D} = \bar{B} \cap \bar{C}$ are isomorphisms $\bar{D} \cong D$, $\bar{D} \cong D'$ commuting with the given one $D \cong D'$. 2) Such a ∇ -commutative system admits a product measure Q.

3) If S' is another ∇ -commutative system verifying the properties of S in (1), then $M(S) \equiv M(S')$.

The proof of 1) and 3) are based on an algebraic theorem concerning biadic algebras that we give next.

We shall write the complete proof of theorem 6 in §7.

THEOREM 7. 1) Given two monadic algebras (B, D, ∇^1) and (C, D', ∇^2) and a fixed isomorphism: D = D' there exists a biadic algebra $(\bar{A}, \bar{B}, \bar{C})$ such that $\bar{B} = B$, $\bar{C} = C$ are isomorphisms which restricted to $\bar{D} = \bar{B} \wedge \bar{C}$ give isomorphisms $\bar{D} = D$, $\bar{D} = D'$ commuting with the given one: D = D'. 2) If there is another biadic algebra $(\bar{A}, \bar{B}, \bar{C})$ with the same proper ties, then the isomorphisms $\bar{B} = \bar{B}$, $\bar{C} = \bar{C}$ obtained through the isomorphisms of \bar{B} , \bar{C} with B and C, have a unique common extension to an isomorphism $\bar{A} = \bar{A}$.

Proof: By identifying the isomorphic algebras D, D' through the given isomorphism we can, without loss of generality, consider only the case that B and C are extensions of the same algebra D. So we have the monadic algebras (B,D,∇^1) and (C,D,∇^2) .

A filter $F
ightharpownewspace{2}{B}$ with the set X_F of ultrafilters of B that contains F, this is a closed subset that represents with the relative topology the quotient algebra B/F (c.f. {5}). In other words, $X_F = S(B/F)$, and in such a way that if \hat{b} is the clopen set that represents b ϵ B, then $\hat{b}
ightharpownewspace{2}{X}_F$ is the clopen set that represents the class in B/F containing b. Given an ultrafilter U in D, let us denote with (U) the filter generated in B by U. Then {X_(U)} is a partition of X. In fact, if m ϵ X corresponds to the ultrafilter M and U = M \cap D, then m ϵ X_(U); if m ϵ X_(U) \cap X_(U'), then M \cap D \supset U, U', which implies U = U'.

Moreover, the monadic operator ∇^1 corresponds with saturation with respect to the partition $\{X_{(U)}\}$; i.e. if $a \in B$, $\nabla^1 a = sat \hat{a} = = \bigcup \{X_{(U)}; \hat{a} \land X_{(U)} \neq \emptyset\}$.

We include the proof of this well-known fact (c.f. {2}) for the sake of completeness. If $d \in D$, then $d \in U$ iff $X_{(U)} \subset \hat{d}$, and it is equivalent to $X_{(U)} \cap \hat{d} \neq \emptyset$, as it is easy to see. In consequence, $\hat{a} \cap X_{(U)} \neq \emptyset$ iff $\nabla^1 a \supset X_{(U)}$, which is also equivalent to $\nabla^1 a \in U$. In fact, $\nabla^1 a \in U$ implies $\hat{a} \cap X_{(U)} \neq \emptyset$, since otherwise $X_{(U)} \subset \hat{C} \hat{a} = \hat{C} \hat{a}$, i.e. (a ε (U), then for some d ε U, d \leq (a, and hence d $\cap \nabla^1 a = 0$, which contradicts the fact that U is a proper filter of D.

Let us observe that given a set X, a partition I of X and a algebra B of subsets of X stable under the sat operation with respect to I, if D is the algebra of saturated sets of B, we have a monadic algebra (B,D,sat.)

Let us suppose now that an extension A of the algebras B, C exists, such that (A,B,C) is a biadic algebra; it is easy to verify that the ∇ operators defined by B, C on A coincide on B, C with the previously given. We call $\nabla = \nabla^1 \nabla^2 = \nabla^2 \nabla^1$ to the monadic operator defined by D.

Since B \oplus C applies homomorphically onto A, preserving the identity mappings on B and C, A = B \oplus C/F for a filter F in B \oplus C. Then the Stone space of A is a closed subset T of S(B \oplus C) = X × Y.

Calling $\{Y_{(U)}\}$ the partition of Y associated to the ultrafilters U of D we have, after elimination of superfluous parentheses:

(*)
$$T = \sum_{U} X_{U} \times Y_{U}$$

In fact, if M ε T corresponds to an ultrafilter \overline{M} of A, and M = (M',M") (M' ε X, M" ε Y), considering the homomorphism mentioned above it follows that M' = B. \overline{M} , M" = C. \overline{M} . Therefore, M'.D = M".D = U is an ultra filter of D. Therefore, M ε X_U × Y_U. Conversely, if (M',M") ε X_U×Y_U, M'.D = M"D = U. To see that (M',M") ε T it suffices to see that there is an ultrafilter \overline{M} of A which is a simultaneous extension of M' and M". It is enough to verify that if b ε M', c ε M" then b $\wedge c \neq 0$. But $\nabla b \wedge \nabla c = \nabla (b \wedge c) \varepsilon$ U and therefore is not zero, which implies $b \wedge c \neq 0$.

Let us consider now the following partitions of T:

1) $\{T_x\}$, $x \in X$, $T_x = (\{x\} \times Y) \cap T$ 2) $\{T_y\}$, $y \in Y$, $T_y = (X \times \{y\}) \cap T$ 3) $\{T_U\}$, U ultrafilter of D, $T_U = X_U \times Y_U$. They define the sat operators corresponding to v^1 , v^2 , v of the biadic algebra (A,B,C). This is immediate for v^1 , v^2 . For v we observe that the monadic algebra (A,D) is represented by T and the partition associated with the ultrafilters U of D, which is precisely $\{X_U \times Y_U\}$ as it was shown in the proof of (*).

If we start with another extension A' of B and C such that (A',B,C)is a biadic algebra we get again the same set T representing S(A')because the second member of (*) depends only on (B,D) and (C,D). Then $A \simeq A'$ by a unique common extension of the identity isomorphisms of B and C. This proves 2) (except for isomorphic identifications). The fact that the definitions of T, the partitions 1), 2), 3), and the clopens corresponding to elements of B and C depend only on (B,D) and (C,D) allow us to construct a biadic algebra of sets (A',B',C') such that $B \approx B'$, $C \approx C'$ and these isomorphisms when restricted to D give an isomorphism $D \approx B' \wedge C' = D'$. That will prove (except for isomorphic identifications) the first part of the theorem.

Let us define A' as the algebra of sets of the form a' = $\hat{a} \cap T$, where \hat{a} is a clopen in S(B) × S(C) = X × Y. The algebras B' = { $(\hat{b} \times Y) \cap T_{b \in B}$, C' = { $(X \times \hat{c}) \cap T_{c \in C}$ generate A'.

We define for a' ε A', $\nabla^1 a'$ as the saturated set with respect to $\{T_x\}$, and $\nabla^2 a'$ that obtained from $\{T_y\}$. We must show that $\nabla^1 a' \varepsilon$ B' to show that (A',B') is monadic, the same for ∇^2 , and that ∇^1 , ∇^2 commute. To this end, let us prove:

6)
$$\nabla^1 c' = (\nabla c \times Y) \wedge T$$
, for c' ε C'.

 $\nabla^1 c' = \nabla^1 (\sum X_U \times (\widehat{c} \cap Y_U))$, where \widehat{c} is the projection of c' on Y, and the star means that the sum is extended to those U such that $\widehat{c} \cap Y_U \neq \emptyset$. Therefore, $\nabla^1 c' = \sum * (\nabla^1 (X_U \times (\widehat{c} \cap Y_U)) = [(\sum * X_U) \times Y] \cap T$. Then, since $\widehat{c} \cap Y_U \neq \emptyset$ is equivalent to $\widehat{\nabla c} = Y_U$, i.e. to $\nabla c \in U$, we get $\sum * X_U = \widehat{\nabla c}$, which proves the formula.

To show that $\nabla^1 a' \in B'$, for every $a' \in A'$, it suffices to see it for $a' = b' \cap c'$, $b' \in B'$, $c' \in C'$; $\nabla^1 (b' \cap c') = b' \cap \nabla^1 c' =$

= $[(\hat{b} \cap \nabla \hat{c}) \times Y] \cap T \in B'$.

If b' = c', then from (6) b' = $(\hat{d}_1 \times Y) \cap T = (X \times \hat{d}_2) \cap T$, which implies $\hat{d}_1 = \hat{d}_2$. Therefore, D' = B' \cap C' is defined as those sets of A' such that are saturated with respect to $\{T_U\}$ and project on the same ele - ment of D. Therefore, (6) means that $\nabla^1 c' \in D'$, which is equivalent to the commutation of ∇^1 , ∇^2 (c.f. $\{1\}$).

AN APPROXIMATION PROCESS.

The generalized product measure Q, when it exists, and in general the finitely additive product measure Q, associated with a ∇ -commutative system S, can be obtained as a limit of simpler measures in the way described in the next theorem. We need the following preliminary result:

PROPOSITION 1. i) If F is a finite part of (A,D,∇) , a monadic algebra, and a_1, \ldots, a_r , the atoms of the Boolean subalgebra generated by F, then a_1, \ldots, a_r , $\nabla a_1, \ldots, \nabla a_r$ generate a subalgebra A_0 which is the

the least one containing F and stable for ∇ , (i.e., $\nabla A_0 \subset A_0$).

ii) Assume A₀ is a finite subalgebra of the biadic algebra (A,B,C, ∇^1, ∇^2), stable for $\nabla = \nabla^1 \nabla^2$, D = B \land C. If B₀ = B \land A₀, C₀ = C \land A₀ then there exist operators ∇^1_0, ∇^2_0 , such that (A₀, B₀, C₀, ∇^1_0, ∇^2_0) is a biadic algebra.

iii) Suppose that $(A_{\lambda})_{\lambda \in \Lambda}$ is the family of all the finite subalgebras of A stable for ∇ , ordered by inclusion, and generated by their elements belonging to B or C. Then:

- 1) every $(A_{\lambda}, B_{\lambda}=A_{\lambda} \cap B$, $C_{\lambda}=A_{\lambda} \cap C$, $\nabla_{\lambda}^{1}, \nabla_{\lambda}^{2}$) is a biadic algebra,
- 2) $(A_{\lambda})_{\lambda \in \Lambda}$ is filtering, if ordered by inclusion,
- 3) $\bigcup_{\lambda} A_{\lambda} = A$.

Proof: i) It is evident that every subalgebra containing F and stable for ∇ must contain the a_i 's and the ∇a_i 's. Then, it suffices to prove that A_o is stable for ∇ . Every atom of A_o is of the form $a = a_i \wedge \bigwedge_j \nabla a_j$ where j runs on some indices 1,...,r. Therefore, $\nabla a = \nabla a_i \wedge \bigwedge_j \nabla a_j \in A_o$.

To finish the proof it suffices to observe that every element of A_o is a union of atoms and that ∇ distributes over the union.

ii) ∇_{0}^{i} exists because A_{0} is finite. $a \in A_{0}$ implies $\nabla a = \nabla^{2} \nabla^{1} a \leq \nabla^{2} \nabla^{0} a \leq \nabla_{0} a = \nabla a$, where ∇_{0} denotes the operator relative to $D \cap A_{0}$. Then $\nabla_{0}^{2} \nabla_{0}^{1} = \nabla_{0}$ and analogously, $\nabla_{0}^{1} \nabla_{0}^{2} = \nabla_{0}$.

iii) It follows from i) and ii) and the observation that every $a \epsilon B \vee C$ belongs to a finite subalgebra generated by a finite set $F \cup G$ with $F \subset B$, $G \subset C$.

Given the ∇ -commutative system $S = (\Omega, A, B, C, P_1, P_2)$ let $M(S) = (A_0, B, C, P_1, P_2)$ be the associated Boolean measure structure. We apply proposition 1 to the biadic algebra (A_0, B, C) to get the filtering family $(A_\lambda)_\lambda$ described in iii). For each λ we select a representative $S_\alpha = (\Omega, A_\lambda, B_\lambda, C_\lambda, P_1, P_2)$ of $(A_\lambda, B_\lambda, C_\lambda, P_1, P_2)$, that means: B_λ, C_λ are finite subalgebras of B, C such that $M(S_\alpha) = (A_\lambda, B_\lambda, C_\lambda, P_1, P_2)$, $A_\lambda = B_\lambda \vee C_\lambda$. Here $\alpha = A_\lambda$, and we assume the α 's ordered by inclusion of the A_λ 's. Ussing theorem 5 we know that the measure Q_α associated to S_α is defined by

$$Q_{\alpha}(B.C) = \sum_{i} P_{1}(BD_{i}) P_{2}(CD_{i})/P(D_{i})$$

where the sum is on the atoms of \mathcal{D}_{λ} (it defines a comarginal commutative measure on S_{α}). Denote with G_{α} the conditional expectation operator relative to \mathcal{D}_{λ} and the probability Q_{α} in the system S_{α} . Then THEOREM 8. In a ∇ -commutative system S, it holds:

i) For $B \in B$, $C \in C$, $G_{\alpha} ^{1}_{B} \longrightarrow Gl_{B}$ and $G_{\alpha} ^{1}_{C} \longrightarrow Gl_{C}$, when $\alpha \uparrow$, uniformely a.e. [P].

116

ii)
$$Q(A) = \lim_{\alpha} Q_{\alpha}(A)$$
, for any $A \in B \vee C$.

Proof: ii) follows from i) and

$$Q_{\alpha}(B.C) = \int G_{\alpha} \mathbf{1}_{B} \cdot G_{\alpha} \mathbf{1}_{C} dP \longrightarrow \int G \mathbf{1}_{B} \cdot G \mathbf{1}_{C} dP = Q(B.C)$$

i) For B ε B_{λ}, α = A_{λ}, we have G_{α}1_B = $\sum_{j} (P_1(BD_j)/P(D_j)) 1_{D_j}$, where the sum is on the atoms of \mathcal{D}_{λ} , since $Q_{\alpha}(BD_j) = \int_{D_j} G_{\alpha} 1_B dP =$

$$= P(BD_j)/P(D_j) \cdot Q_{\alpha}(D_j) = P_1(BD_j) \quad (P_1 = Q_{\alpha} \text{ on } B_{\lambda}).$$

Given G1_B, let us divide the set of real numbers on intervals $I_i = [m_i, M_i)$ of length $\varepsilon > 0$. Only for a finite number of them $D_i = (G1_B)^{-1}(I_i) \neq \emptyset$ [P]. Consider any finite system S_α such that $D_i \varepsilon \mathcal{D}_\lambda$ for all those D_i . Call $\{D_{ij}\}$ the family of atoms of \mathcal{D}_λ contained in D_i .

Since,
$$m_i P(D_{ij}) \leq \int_{D_{ij}} G \mathbf{1}_B dP \leq M_i P(D_{ij})$$

we obtain: $m_i P(D_{ij}) \leq P_1(BD_{ij}) \leq M_i P(D_{ij})$. Therefore on D_{ij} : |G 1_B - $P_1(BD_{ij})/P(D_{ij})| \leq \varepsilon$, a.e [P]. Then

$$|\sum_{i}\sum_{j} P_{1}(BD_{ij})/P(D_{ij}) |_{D_{ij}} - G |_{B}| \le \varepsilon$$
 a.e [P] QED

REMARK: If a comarginal probability P exists on A we have $G_{\alpha} = 1_{B} = E_{p}(1_{B} | \mathcal{D}_{\lambda})$. Using a result from martingale theory due to Helms {3}, we know that the G_{α} 's form a uniformly integrable martingale converging in L^{1} to $G_{1_{B}}$, which implies ii).

7. PRODUCT MEASURE

ror special cases of ∇ -commutative systems we can assert that a product measure exists.

One of these cases, the discrete case, was considered in theorem 5. Using remark 4) of §2 we have also:

THEOREM 9. If the ∇ -commutative system S = (Ω , A, B, C, P₁, P₂) is such

that there are semi-compact classes $K_B \subset B$, $K_C \subset C$ with the property of approximation and $K_B \cdot K_C = \{K.L ; K \in K_B, L \in K_C\}$ is also semi-compact, then Q(B.C) = $\int Gl_B \cdot Gl_C dP$ defines, when extended to A, a probality, i.e. the product measure on S.

Another case in which we can assert the existence of product measure is referred in theorem 6, the proof of which we give now:

Proof of theorem 6. We will suppose, like in the proof of theorem 7, that D is identified to D' through the given isomorphism. By theorem 7 we have an extension algebra A_0 of B, C such that (A_0, B, C) is biadic and $B \cap C = D$. Then we have in the Stone space T of A_0 the algebras of clopens \hat{A}_0 , \hat{B} , \hat{C} relative to A_0 , B, C. We set $B = \tau(\hat{B})$, $C = \tau(\hat{C})$ and $A = \tau(\hat{A}_0) = \tau(B, C)$. We define in \hat{B} , \hat{C} the measures P_1 , P_2 given in B, C in the canonical way and extend them to B, C.

Let us prove that $S = (T, A, B, C, P_1, P_2)$ is a ∇ -commutative system. If $\hat{b} \in \hat{B}$ and $\hat{c} \in \hat{C}$ are such that $\hat{b} \cdot \hat{c} = \emptyset$, i.e. $b \wedge c = 0$ in A; we have $0 = \nabla^2(b \wedge c) = \nabla^2 b \wedge c = \nabla^2 \nabla^1 b \wedge c = \nabla b \wedge c$. Then $\nabla \hat{b} \cdot \hat{c} = \emptyset$, where $\nabla \hat{b} \in \hat{E} \wedge \hat{C}$ contains \hat{b} . Hence $\{G1_{\hat{b}}^{*} > 0\} \subset \nabla \hat{b}$ [P], and then

 $\int G1_{\hat{\mathbf{b}}} \cdot G1_{\hat{\mathbf{c}}} \, dP = \int G1_{\hat{\mathbf{b}}} \cdot 1_{\hat{\mathbf{c}}} \, dP_2 \leq P_2(\hat{\mathbf{vb}}, \hat{\mathbf{c}}) = 0$

where G is defined on B(C)-measurable functions is the expectation operator relative on $\mathcal{D} = B \cap C$.

But $\hat{B} \subset B$ and $\hat{C} \subset C$ have the approximation property. Using remark 3 of §2 we can assert that $Q(B.C) = \int G1_B.G1_C dP$ is a finitely additive comarginal measure on S, and from theorem 2, S is a ∇ -commutative system.

Q being finitely additive on $B \vee C \supseteq \hat{A}_{o}$ is a fortiori σ -additive on the algebra of clopens \hat{A}_{o} and then can be extended uniquely to $A = \tau(\hat{A}_{o})$.

This proves 1) and 2) of the theorem (except for identifications) 3) follows immediately from theorem 7.

Finally we shall prove:

THEOREM 10. If $S = (\Omega, A, B, C, P_1, P_2)$ is a ∇ -commutative system with the property:

(P) Ω is the only set of B containing a set C ε C with $P_2(C) > 0$, then $Q(B.C) = P_1(B).P_2(C)$ can be extended to a probability on A. ((P) implies that S is simple). In order to prove the σ -additivity of $Q = P_1 \cdot P_2$ it is enough to prove that $\sum_{\alpha} B_{\alpha} \cdot C_{\alpha} = \Omega$ implies $\sum_{\alpha} P_1(B_{\alpha}) \cdot P_2(C_{\alpha}) = 1$. For each pair $\alpha \neq \beta$ of indices the set $N_{\alpha,\beta} = \{x ; (B_{\alpha}C_{\alpha})_x \cdot (B_{\beta}C_{\beta})_x \neq \emptyset \ [P_2]\}$ is contained in a set of B of P_1 -measure zero. In fact, from $B_{\alpha}C_{\alpha}B_{\beta}C_{\beta}=\emptyset$ we have either $P_1(B_{\alpha} \cdot B_{\beta}) = 0$ or $P_2(C_{\alpha} \cdot C_{\beta}) = 0$. In the first case, it follows from $N_{\alpha\beta} \subset B_{\alpha}B_{\beta}$, in the second one, since $(B_{\alpha}C_{\alpha})_x \cdot (B_{\beta}C_{\beta})_x$ $\subset C_{\alpha} \cdot C_{\beta}$ for every x, we have $N_{\alpha\beta} = \emptyset$.

Then except for a set B_o of B of P₁-measure zero (> $\bigcup_{\alpha \neq \beta} N_{\alpha\beta}$)

(7)
$$P_2(\bigcup_{\alpha} (B_{\alpha}, C_{\alpha})_x) = \sum_{\alpha} P_2((B_{\alpha}, C_{\alpha})_x) = \sum_{\alpha} P_2(C_{\alpha}) \cdot 1_{B_{\alpha}}(x)$$
.

Let b_x be the set of the partition of Ω defined by the sets B_α such that b_x 3 x

(8)
$$b_{\mathbf{x}} \subset \bigcup_{\alpha} (B_{\alpha}, C_{\alpha})_{\mathbf{x}}$$

In fact, let $y \in b_x$ and suppose $y \in B_\alpha \cdot C_\alpha$, then $B_\alpha \supset b_x \ni x$ and this implies $(B_\alpha \cdot C_\alpha)_x = C_\alpha$. Since $y \in C_\alpha$ we have $y \in \bigcup_\alpha (B_\alpha \cdot C_\alpha)_x$. This proves (8).

From $B \ni b_x \subset \bigcup_{\alpha} (B_{\alpha}.C_{\alpha})_x \in C$, and the assumed property (P) we have $P_2(\bigcup_{\alpha} (B_{\alpha}.C_{\alpha})_x) = 1$ for every x.

Hence, by virtue of (7) we have $\forall x \notin B_0$: $\sum_{\alpha} P_2(C_{\alpha}) 1_{B_{\alpha}}(x) = 1$. i.e. $\sum_{\alpha} P_2(C_{\alpha}) . 1_{B_{\alpha}} = 1_{\Omega} [P_1]$. By integration with respect to P_1 :

 $\sum_{\alpha} P_2(C_{\alpha}) \cdot P_1(B_{\alpha}) = 1 \qquad \text{QED}.$

Remark: Condition (P) is equivalent to: (P*) if $C \in C$ contains $B \neq \emptyset$, $B \in B$, then $C = \Omega$ a.e. $[P_2]$.

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SÔBRE O PÔSTO DE UM MÓDULO (1)

Jorge Aragona e Artibano Micali

O objetivo desta nota é o de generalizar um resultado sôbre formas lineares que se encontra em {2} . No que segue, todo anel é supo<u>s</u> to comutativo e com elemento unidade.

1. PRELIMINARES.

Sejam A um anel e M um A-módulo. Designaremos por M* o dual de M e, se A fôr un anel de integridade, por t(M) o sub-A-módulo de tor ção de M. É conhecido que t(M) é o núcleo da aplicação natural M \longrightarrow M \mathfrak{S}_A K, onde K é o corpo de frações de A.

Seja A \longrightarrow K um homomorfismo de anéis de A num corpo K que transforma elemento unidade em elemento unidade.Então K pode ser munido, de uma maneira evidente, de uma estrutura de A-módulo. Para todo A-módulo M, o K-pôsto de M é definido como sendo a dimensão do Kespaço vectorial M \mathfrak{S}_A K e notaremos $r_K(M) = [M \mathfrak{S}_A K : K]$. Se A fôr um anel de integridade en K o corpo de frações de A, falaremos simplesmente do pôsto de M e indicaremos com $r(M) = [M \mathfrak{S}_A K : K]$.

LEMA 1. Sejam A um anel $e \ 0 \longrightarrow M' \longrightarrow M \longrightarrow M'' \longrightarrow 0$ uma sequè<u>n</u> cia exata de A-módulos. Temos $r_{K}(M) \le r_{K}(M') + r_{K}(M'') e$ se $Tor_{1}^{A}(M'',K) = 0$, então $r_{K}(M) = r_{K}(M') + r_{K}(M'')$.

Com efeito, é suficiente ver que se tem a sequência exata de K-espaços vectoriais

 $\dots \longrightarrow \operatorname{Tor}_{1}^{\mathbb{A}}(\mathsf{M}'',\mathsf{K}) \longrightarrow \mathsf{M}' \, \boldsymbol{\mathfrak{S}}_{\mathbb{A}} \, \mathsf{K} \longrightarrow \mathsf{M} \, \boldsymbol{\mathfrak{S}}_{\mathbb{A}} \, \mathsf{K} \longrightarrow \mathsf{M}'' \, \boldsymbol{\mathfrak{S}}_{\mathbb{A}} \, \mathsf{K} \longrightarrow 0$

COROLÁRIO. Se A fôr um anel de integridade e $0 \longrightarrow M' \longrightarrow M \longrightarrow M'' \longrightarrow 0$ uma sequência exata de A-módulos, então r(M) = r(M') + r(M'').

Isto resulta do fato de que K, corpo de frações de A, é um A-módulo plano.

Observemos finalmente que se A é um anel de integridade e $a \neq 0$ um ideal de A, então r(a) = 1. Com efeito, é claro que r(a) \leq 1. Supondo que r(a) = 0, deduziríamos que a \mathfrak{S}_{A} K = 0, onde K é o corpo

(1) Trabalho realizado com auxílio de FAPESP Proc. Matemática 66/038.

de frações de A, e como a está contido em a Θ_A K, seguiria que a=0.

2. MODULOS DE POSTO n.

TEOREMA. Sejam A um anel de integridade, M um A-módulo e $n \ge 1$ um inteiro. As seguintes condições são equivalentes:

- i) existe uma família livre $(x_i)_{1 \le j \le n}$ de elementos de M tal que Ann (M/($x_1, ..., x_n$)A) $\ne 0$;
- ii) existe uma família (f_i)_{1≤i≤n} de formas linares sôbre M e uma família (x_i)_{1≤i≤n} de elementos de M tais que f₁(x₁) = ... = f_n(x_n) ≠ 0 , f_i(x_i) = 0 se i ≠ j e

 $f_{i}(x_{i})y = \sum_{j=1}^{n} f_{j}(y)x_{j} \text{ para todo } y \in M \text{ e para } i = 1, \dots, n ;$ iii) r(M) = n;

iv) existe uma família $(f_i)_{1 \le i \le n}$ de formas linares sôbre M, duas a duas distintas, tais que: $t(M) = \bigcap_{i=1}^{n} Ker(f_i) e$, para todo j, $1 \le j \le n$, $\bigcap_{i \ne i} Ker(f_i) \ne t(M)$

(i) \implies (ii). Como Ann(M/(x₁,...,x_n)A) \neq 0, existe um elemento a ϵ A, a \neq 0 tal que aM \subset (x₁,...,x_n)A. Logo, para todo y ϵ M, se tem ay = $\sum_{i=1}^{n} f_i(y) x_i$ onde $f_i(y) \epsilon$ A para todo i e como a familia (x_i)_{1 \leq i \leq n} ϵ livre, $f_i \epsilon$ uma aplicação linear de M em A para todo i.Com efeito, (a - $f_i(x_i) x_i - \sum_{j \neq i} f_j(x_i) x_j = 0$ implica que $f_i(x_i) = a \neq 0$ para todo i e $f_j(x_i) = 0$ para i \neq j. Temos assim $f_i(x_i) y = \sum_{j=1}^{n} f_j(y) x_j$ para todo y em M.

De outro lado, tomando $y_1 e y_2 em M$, se tem $\sum_{i=1}^{n} f_i (y_1 - y_2) x_i =$ = $f_i (x_i) (y_1 + y_2) = f_i (x_i) y_1 + f_i (x_i) y_2 = \sum_{i=1}^{n} f_i (y_1) x_i +$

+ $\sum_{i=1}^{n} f_i(y_2) x_i$, logo $f_i(y_1 + y_2) = f_i(y_1) + f_i(y_2)$ para todo i. <u>A</u> nàlogamente, se c ε A e y ε M, deduzimos que $f_i(cy) = c.f_i(y)$ para todo i.

(ii) \implies (i) Seja $\sum_{j=1}^{n} a_{j} x_{j} = 0$ onde os a_{j} estão em A. Resulta, da igualdade precedente, que $0 = f_{i}(\sum_{j=1}^{n} a_{j} x_{j}) = a_{i} f_{i}(x_{i})$ e como $f_{i}(x_{i}) \neq 0$ e A é um anel de integridade, então $a_{i} = 0$ para todo i. Supondo que Ann(M/(x_{1},...,x_{n})A) = 0 e reduzindo a relação $f_{i}(x_{i})y=$
$$\begin{split} &\sum_{j=1}^{n} f_{j}(y) x_{j} \text{ para todo } y \in M, \text{ modulo } (x_{1}, \dots, x_{n}) A, \text{ deduziriamos} \\ &f_{i}(x_{i}) \in \operatorname{Ann}(M/(x_{1}, \dots, x_{n}) A) \text{ de onde, } f_{i}(x_{i}) = 0 \text{ para todo } i. \end{split}$$
 $(i) \implies (iii). \text{ Como } A \notin \text{ um anel de integridade } \acute{e} a \text{ familia } (x_{i})_{1 \leq i \leq n} \\ \acute{e} \text{ livre, a sequência exata de } A-modulos } 0 \longrightarrow (x_{1}, \dots, x_{n}) A \longrightarrow M \rightarrow \\ \rightarrow M/(x_{1}, \dots, x_{n}) A \longrightarrow 0 \text{ nos da } r(M) = r((x_{1}, \dots, x_{n}) A) + r(M/(x_{1}, \dots, x_{n}) A) = \\ n + r(M/(x_{1}, \dots, x_{n}) A). Além \text{ disso como } Ann(M/(x_{1}, \dots, x_{n}) A) \neq 0 , \\ então (M/(x_{1}, \dots, x_{n}) A) \otimes_{A} K = t(M/(x_{1}, \dots, x_{n}) A) \otimes_{A} K = 0. \text{ Logo } , \\ r(M) = n. \end{split}$

(iii) \implies (i). Se r(M) = n, é claro que o K-espaço vectorial M \bigotimes_A K tem sempre uma base do tipo $(x_i \otimes 1)_{1 \le i \le n}$ e isto implica que a família $(x_i)_{1 \le i \le n}$ é livre. De outro lado, para todo y ε M, y \neq 0, pode mos escrever y \otimes 1 = $\sum_{i=1}^{n} (a_i/b_i)(x_i \otimes 1)$ onde os a_i/b_i estão em K. Pondo b = $\prod_{i=1}^{n} b_i$ e $c_i = a_i \prod_{i \ne j} b_j$, a relação precedente é equiva - lente à relação (by - $\sum_{i=1}^{n} c_i x_i$) \otimes 1 = 0. Isto implica (cf. {1}, cap 2, prop. 4) que by - $\sum_{i=1}^{n} c_i x_i$ não é livre em M, logo que existe um elemento c ε A, c \neq 0 tal que c(by - $\sum_{i=1}^{n} c_i x_i$) = 0. Logo, cb ε Ann(M/(x_1, \ldots, x_n)A) e cb \neq 0, uma vez que A é um anel de integridade.

(iv) \Longrightarrow (iii). Com efeito, para todo j, $1 \le j \le n$, existe um elemento $x'_{j} \in \bigcap_{i \ne j} \operatorname{Ker}(f_{i}) - t(M)$ tal que $x'_{j} \notin \operatorname{Ker}(f_{j})$. As hipóteses feitas implicam que a família $(x'_{j})_{1 \le j \le n} \notin t$ livre. Pondo $x_{i} = f_{1}(x'_{1}) \dots f_{i}(x'_{i}) \dots f_{n}(x'_{n})x'_{i}$ para todo i, segue-se que a famí lia $(x_{i})_{1 \le i \le n} \notin t$ livre, que $f_{j}(x_{i}) = 0$ se i $\ne j$ e que $f_{1}(x_{1}) = f_{2}(x_{2}) \dots$ $= f_{n}(x_{n}) \ne 0$. Seja f: $M \longrightarrow \bigoplus_{i=1}^{n} f_{i}(M)$ a aplicação A-linear definida por $f(x) = (f_{i}(x))_{1 \le i \le n}$ para todo x em M e J = Im(f). A sequência exata de A-módulos $0 \longrightarrow t(M) \longrightarrow M \longrightarrow J \longrightarrow 0$ nos dá r(M) = r(J). Consideremos a família $(e_{i})_{1 \le i \le n}$ de elementos de J definida por $e_{j} =$ $(f_{i}(x_{j}))_{1 \le i \le n}$. Mostremos que a família $(e_{i})_{1 \le i \le n}$, onde os c_{i} estão em A, implica que $c_{i}f_{i}(x_{i}) = 0$, logo que $c_{i} = 0$ para todo i. De outro lado, se y ε J existe un elemento x ε M tal que $y=(f_{i}(x))_{1 \le i \le n}$ Se tomarmos $c = -f_{1}(x_{1}) = -f_{2}(x_{2}) = \dots = -f_{n}(x_{n}) \ne 0$, temos a r<u>e</u> 122

lação não trivial cy + $\sum_{i=1}^{n} f_i(x) e_i = 0$ entre os $e_i e_j$. Isto nos mostra que a família $(e_i)_{1 \le i \le n}$ é livre maximal, logo que r(J) = n.

(ii) \implies (iv). Trivial.

OBSERVAÇÕES.

(1) Sejam A um anel de integridade, M um A-módulo e $n \ge 1$ um inte<u>i</u> ro. Se uma qualquer das condições equivalentes do teorema precede<u>n</u> te fôr verificada, então M não pode ser um módulo de torção. Com <u>e</u> feito, se M = t(M), então r(M) = 0, o que é absurdo, uma vez que r(M) = $n \ge 1$.

(2) No caso n = 1, $f_1 = f$, em (iv) \implies (iii) se tem a sequência exa ta de A-módulos $0 \longrightarrow t(M) \longrightarrow M \longrightarrow f(M) \longrightarrow 0$ e portanto, as condições $t(M) \neq M$ e f $\neq 0$ são equivalentes.

3. A x-CONDIÇÃO

1. Em seguida, vamos mostrar como o teorema acima se relaciona com a noção de x-condição introduzida em {2}. Sejam A um anel de integridade, M um A-módulo e f uma forma linear sôbre M. Se existir um x em M tal que $f(x) \neq 0$ e f(x)y = f(y)x para todo y ε M, diremos que a forma linear f obedece à x-condição.

EXEMPLOS.

1) Se f é uma forma linear injectiva sôbre M, então f obedece à xcondição, com x ε M - {0}, x qualquer. Reciprocamente a x-condição não implica que f seja injectiva, mas sòmente que Ker(ψ) = 0:f(x)A. Vemos assin que se 0:f(x)A = 0 as duas noções coincidem.

2) Daremos aquí um exemplo de forma linear que obedece à x-condição mas que não é injectiva. Para isto, sejam A = Z o anel dos inteiros racionais, M = Z × Z/(2) (produto direto) e f a forma linear sôbre M definida por f(m, \bar{n}) = 2m para todo (m, \bar{n}) ϵ M. É claro que f é uma linear sôbre M e se consideramos o elemento x = (1, $\bar{0}$) ϵ M,f(x) $\neq 0$. De outro lado, para todo y = (m, \bar{n}) ϵ M se tem f(x)y = 2(m, \bar{n}) = = (2m, $\bar{0}$) e f(y)x = 2m(1, $\bar{0}$) = (2m, $\bar{0}$). Assim f obedece a (1, $\bar{0}$)-cond<u>i</u> ção. O núcleo de f é o sub-A-módulo de M formado pelos pares (0, $\bar{0}$) e (0, $\bar{1}$).

3) Consideremos o anel A = Z, o A-módulo M = Z × Z (produto direto) e a forma linear f sôbre M definida por f(m,n) = m + n, para todo (m,n) ε M. Suponhamos que exista um x = (p,q) ε M tal que x obedeça à x-condição. Para todo y = (m,n) ε M se deve ter (p+q)(m,n)= = (m+n)(p,q), isto é np = mq. Ora, uma tal relação é impossível para todo $(m,n) \in M$, exceto se p = q = 0, o que daria x = 0 e portanto f(x) = 0.

2. Indíquemos com $M^*(x)$ o sub-A-módulo de M^* das formas lineares que obedecem à x-condição. É lógico que podemos excluir sempre o caso x=0, uma vez que $M^*(0) = \emptyset$. Obtem-se assim:

$$\bigcup_{x \in M} M^*(x) \subset M^*$$
, logo $\langle \bigcup_{x \in M} M^*(x) \rangle \subset M^*$

onde indicamos com <S> , se S é um sub-conjunto de M, o sub-A-módulo de M gerado por S. Veremos, no N° 4, que em geral:

$$\langle \bigcup_{\mathbf{x} \in \mathbf{M}} M^*(\mathbf{x}) \rangle \subseteq M^*$$

isto é, não é possível "aproximar" toda forma linear por formas $1\underline{i}$ neares que obedecem a x-condições, x em M.

3. Do teorema precedente, resulta uma caracterização dos módulos munidos de uma forma linear que satisfaz a uma x-condição.

COROLÁRIO. Sejam A um anel de integridade e M um A-módulo. As co<u>n</u> dições seguintes são equivalentes:

(1) existe um elemento x em M tal que $Ann({x}) = 0 e Ann(M/Ax) \neq 0$.

(2) existe uma forma linear f sôbre M que obedece à x-condição.

- (3) r(M) = 1.
- M não é um módulo de torção e existe uma forma linear f sôbre M tal que Ker(f) = t(M).

4. CONTRA-EXEMPLO A FÓRMULA $\langle \bigcup_{x \in M} M^*(x) \rangle = M^*$.

Suponhamos que a fórmula acima seja verdadeira. Toda forma linear f $\in M^*$ se escreve $f = \sum_i f_i$ (soma finita), onde $f_i \in M^*(x_i)$ para to to i. Logo $f_i(x_i)y = f_i(y)x_i$ para todo y $\in M$. Isto nos mostra que $f_j(x_i)y - f_j(y)x_i \in Ker(f_i) = t(M)$ (pelo corolário) isto é, existe um elemento c $\in A$, c $\neq 0$ tal que $f_j(cx_i)y = f_j(y)(cx_i)$ para todo y. Logo $f_j \in M^*(cx_i)$ e portanto $f \in M^*(cx_i)$. Mas, o exemplo 3) do N°3 nos mostra que, em geral, isto não é possível.

NOTA: La traducción de "pôsto" es "rango".

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Este volumen constituye una excelente y apropiada exposición de las aplicaciones de procesos estocásticos a la biología y la medicina . La presentación de los temas se encuentra completamente al día.

La organización en general se divide en tres largos capítulos. Los dos primeros tratan la teoría de procesos estocásticos de parámetro discreto y continuo respectivamente los cuales constituyen una in troducción al último capítulo dedicado a las aplicaciones, que son muy variadas. Como dato ilustrativo cabe citar a modo de ejemplo los estudios de los modelos de población, migraciones epidemia, modelo de Karlin-McGregor en los procesos estocásticos de la evolu ción, modelo estocástico de la contracción muscular y hasta un muy reciente modelo de la carie dentaria debida a Lu (1966).

En síntesis, es una obra muy completa y muy competente en el tema , tal vez una de las primeras en su tipo. El lector se enfrenta con una larga bibliografía.

E. Marchi.

THEORY OF RANDOM FUNCTIONS, por A.Blanc-Lapierre y R.Fortet, Vol.I, Gor don and Breach, New York-London-Paris, 1967, 454 pgs., 29, 50 dólares.

Este libro es una traducción al inglés por J. Gani de la primera m<u>i</u> tad del bien conocido texto "Théorie des fonctions aléatoires". Además de un apéndice sobre las nociones matemáticas básicas necesarias para desarrollar el material expuesto contiene: Cap. I. Intr<u>o</u> ducción práctica al estudio de las funciones aleatorias. Cap. II. Axiomas, conceptos básicos y teoremas fundamentales de la teoría de la probabilidad. Cap. III. Introducción general a las funciones <u>a</u> leatorias. Cap. IV. Introducción general a los procesos estocást<u>i</u> cos: funciones aleatorias con incrementos independientes. Cap. V. Funciones aleatorias derivadas de procesos de Poisson. Cap. VI. Procesos de Markov. Cap. VII. Cadenas de Markov. Funcionales ad<u>i</u> tivas de un proceso de Markov.

La encuadernación y la impresión son excelentes y el volumen II-que completaría la traducción del original francés - es prometido.

R. Panzone.







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Los volúmenes III, IV, V y VI comprenden los siguientes fascículos separados.

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Además han aparecido tres cuadernos de Miscelánea Matemática.

INDICE

Volumen 24, N° 2, 1968

Induced sheaves	and Grothendieck topologies	
	por J. J. Martínez	67
Note on Galois	extension over the center	
	por Manabu Harada	91

Volumen 24, N° 3, 1969

A mean value theorem and Darboux's property for the derivative	
of an additive set function with respect to a measure on E^{H}	
por R. J. Easton y S. G. Wayment	97
Notes on comarginal probability measures	÷
por A. Diego y R. Panzone	103
Sôbre o pôsto de um módulo	
por J. Aragona y A. Micali	119
Comentarios Bibliográficos	125