# a Generalization of the riemann-roch theorem 

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## 1. Introduction

In this paper ${ }^{1}$ a Riemann-Roch theorem is proved for a module, over a function field $K$, which is under the action of simple algebras over $K$. Specialization of this module leads on one hand to the Riemann-Roch theorem of E. Witt [16] for simple algebras over $K$, and on the other hand to an extension of A. Weil's Riemann-Roch theorem for matrices over function fields [15], in the case that his "signature" is taken to be identically 1 . In each case the constant field is allowed to be arbitrary.

There is also a brief account (in §2), partly new in method, of the arithmetic of simple algebras over $K$. In §3 our generalization of the Riemann-Roch theorem is proved for a certain module over the function field $K$. In $\S 4$ this module is taken to be a simple algebra $A$ over $K$; a restriction of the definition of divisor then leads to a suitably specific form of the Riemann-Roch theorem for $A$. Related questions-the different, the Riemann-Hurwitz formula, and a genus-like invariant of $A$-are then discussed. Finally, in $\S 5$, it is shown that our Riemann-Roch theorem for $A$ implies that of Witt [16]. The paper concludes with a theorem extending the generalized Riemann-Roch theorem of Weil [15] (when his "signature" is trivial) for matrices over function fields.

Part of the origin of this kind of investigation is in the papers of Hecke [7, 8], Chevalley and Weil [3], and Weil [14], which are concerned with the problem of decomposing into its irreducible parts a certain natural representation of $G / H(N)$, where $G$ is the modular group and $H(N)$ the subgroup of matrices congruent $(\bmod N)$ to the identity, as linear transformations of the space of "cusp forms" of type $(2, N)$. Since there is a natural isomorphism between this space of cusp forms and the differentials of the first kind of the associated function field $K_{H(N)}$, the problem can be transformed to one in terms of matrices over $K_{H(N)}$.

The methods used here are those of linear topology and duality, first applied to this kind of problem by K. Iwasawa in [10] and particularly [9]. The proofs in §3 are direct generalizations of the proofs of Iwasawa for the corresponding theorems about $K$. Indeed, much of this paper may be thought of as the tensor product of the appropriate spaces over $K$ with [9].

I wish to thank Professor Iwasawa for suggesting to me the problems dealt with here. As I have indicated, his works [9] and [10] made these problems

[^0]easy to solve. I have also profited from several discussions with Professor Iwasawa and with my friend Peter Schweitzer, from whom I have received valuable advice and encouragement. Finally, I wish to thank the referee for several very helpful suggestions.

For the preliminaries, let $K$ denote a function field ${ }^{2}$ with field of constants $F$. Let $\mathfrak{M}$ denote the set of prime divisors of $K, \nu_{P}$ the normalized exponential valuation belonging to $P \in \mathfrak{M}, K_{P}$ the completion of $K$ with respect to $P$, $\mathfrak{o}_{P} \subset K_{P}$ the ring of local integers, $\mathfrak{p}_{P}$ the local prime ideal, $\Omega_{P}=\mathfrak{o}_{P} / \mathfrak{p}_{P}$ the residue-class-field, and $n_{P}=\left[\Omega_{P}: F\right]$ the local degree.

The ring $R$ of valuation vectors of $K$ is defined as the weak direct sum of the $K_{P}$ in the sense that for almost all $P$, the component $a_{P}$ of the valuation vector $a$ must belong to $\mathfrak{o}_{P}$. In $R$ we define the subring $\mathfrak{o}$ as the direct sum of the $\mathfrak{o}_{P}$. We then take the set of all $a \mathfrak{0}$, where $a$ runs through the regular elements of $R$, as a fundamental system of neighborhoods of 0 in $R$, defining thereby a linear topology ${ }^{3}$ on $R$, under which $R$ even becomes a topological ring. $\quad K$ is a discrete subfield of $R$; and there exists an open linearly compact subspace $W$ such that $R=K+W$, a topological direct sum (cf. [9, 10]). A character of $R$ is a continuous, $F$-linear mapping of $R$ into $F, F$ having the discrete topology. The space $X(R)$ of all such characters, with the (linearly) compact-open topology, is a linearly topologized vector space over $F$. A nontrivial character of $R$ which vanishes on $K$ will be called an admissible character of $R$. If $\chi$ is such a character, then a fundamental result of $[9,10]$ is that $R$ is self-dual under the pairing $(a, b)=\chi(a b), a, b \in R$.

We shall so often need the following result from the theory of linear topologies that we state it here as

Lemma 1.1. A linearly topologized space is finite-dimensional if and only if it is linearly compact and discrete.

Let $S_{0}$ be a skew-field of finite rank over the center $K$. For each prime divisor $P$ of $K$, define $S_{0 P}$ to be the tensor product $K_{P} \otimes S_{0}$ over $K$. Then, although $S_{0 P}$ is not always a skew-field, it is a normal simple algebra with $K_{P}$ as center. ${ }^{4} \quad$ As such, $S_{0 P}$ is isomorphic to the full $\mu_{P} \times \mu_{P}$ matrix algebra over some skew-field $S_{P}$ with $K_{P}$ as center; and $S_{P}$ is uniquely determined up to a $K_{P}$-isomorphism.

The valuation $\nu_{P}$ of $K_{P}$ can be extended uniquely to $S_{P}$, via, for instance, the norm of the regular representation of $S_{P} / K_{P}$. We denote by $v_{P}$ the

[^1]uniquely determined normalized exponential valuation of $S_{P}$ which results from this extension. $S_{P}$ is complete with respect to $v_{P}$. We have the result that $\left[S_{P}: K_{P}\right]=e_{P} f_{P}$, where $e_{P}$ is the local ramification index and $f_{P}$ the local rank of $S_{P} / K_{P}$. We define $\mathcal{O}_{P}$ to be the ring of elements in $S_{P}$ with nonnegative valuation, $\mathfrak{B}_{P}$ to be the maximal ideal of non-units in $\mathfrak{O}_{P}$, and $\mathfrak{S}_{P}=$ $\mathcal{O}_{P} / \Re_{P}$, the local skew-field at $P . \mathfrak{S}_{P}$ is of finite rank $f_{P}$ over $\Omega_{P}$.

It will also be useful to recall the following well-known result from the theory of valuations: Let $w_{1}, \cdots, w_{f} \in \mathcal{O}_{P}$ be representatives of a basis of $\Im_{P} / \Omega_{P}$, and let $t \in S_{P}$ be a prime element for $P$, i.e., $v_{P}(t)=1$. Then we have

Lemma 1.2. ${ }^{5}$ The $e_{P} f_{P}$ elements $w_{i} t^{j}$, with $1 \leqq i \leqq f_{P}$ and $0 \leqq j \leqq e_{P}-1$, form a basis of $S_{P} / K_{P}$; and if $a \in S_{P}$ is written $a=\sum a_{i j} w_{i} t^{j}$, with $a_{i j} \in K_{P}$, then, for every integer $m, v_{P}(a) \geqq m e_{P}$ if and only if all $\nu_{P}\left(a_{i j}\right) \geqq m$.

Lemma 1.3. If the skew-field $S_{P}$ has rank $m_{P}^{\prime 2}$ over the center $K_{P}$, then $e_{P} \mid m_{P}^{\prime}$ and $m_{P}^{\prime} \mid f_{P},[18]$.

Proof. If $t$ is a prime element for $P$, then $K_{P}(t) / K_{P}$ has ramification index $e_{P}$, which divides its degree, which in turn divides $m_{P}^{\prime}$ [1, p. 53]. Since $e_{P} \mid m_{P}^{\prime}$, we have $m_{P}^{\prime 2}=e_{P} f_{P}$ dividing $m_{P}^{\prime} f_{P}$, or $m_{P}^{\prime} \mid f_{P}$.

Let $A$ be an algebra over the center $k$. The reduced trace $T$ of $A / k$ is always nontrivial if $A$ is semisimple (in fact, the discriminant is nonzero). ${ }^{6}$ In particular, $T$ is a $k$-linear mapping of $A$ to $k$ such that $T(a b)=T(b a)$ for all $a, b \in A$.

## 2. The arithmetic of simple algebras over function fields

In this section we present a brief account of the maximal orders ${ }^{7}$ and ideals of a normal simple algebra $A$ over $K$. The results are known (cf. [5] and [12]), but some of the present proofs are simpler than the older ones.

We first define the ring $\widetilde{A}$ of valuation vectors of $A$ as the tensor product $\tilde{A}=R \otimes A$ of $R$ and $A$ over $K$. If $u_{1}, \cdots, u_{n}$ is a basis of $A / K$, then $\widetilde{A}=$ $R \otimes u_{1}+\cdots+R \otimes u_{n}$; thus we may give $\widetilde{A}$ the linear topology of a direct sum of copies of $R$. This topology makes $\widetilde{A}$ a topological ring and is independent of the basis of $A / K$ chosen above.

Let $x$ be an element of $K$ not in $F$. Call a prime divisor $P$ of $K$ finite if $\nu_{P}(x) \geqq 0$. An order of $A$ relative to $F[x]$ is a subring of $A$, finitely generated over $F[x]$, containing $F[x]$, and spanning $A / F(x)$. In this section $Q$ will always denote that prime divisor of $F(x)$ obtained by projecting to $F(x)$ that $P \in \mathfrak{M}$ which is mentioned in the same context as $Q . \quad R_{0}$ is defined as the ring of valuation vectors of $F(x), \mathfrak{o}_{0}$ the ring of integers of $R_{0}$, and $\mathfrak{o}_{Q}$ the local integers of the completion $F(x)_{Q}$. One of the basic results of the

[^2]classical theory is that $R$ is topologically isomorphic to the tensor product $R_{0} \otimes K$ over $F(x)$, the latter having the direct-sum topology; we write ${ }^{8}$
\[

$$
\begin{equation*}
R=R_{0} \otimes_{F(x)} K \tag{2.1}
\end{equation*}
$$

\]

We now prove, letting $\tilde{A}_{P}$ denote $K_{P} \otimes A$, the $P$-component of $\tilde{A}$,
Theorem 2.1. The maximal orders of $\tilde{A}_{P}$ relative to $\mathfrak{o}_{P}$ are the same as those relative to $\mathrm{o}_{Q}$ and consist of all the maximal open linearly compact subrings of $\tilde{A}_{P}$.

Proof. Let $L$ be a maximal open linearly compact subring of $\widetilde{A}_{P}$. Then $L+\mathfrak{o}_{P}+\mathfrak{o}_{P} L$ is an open linearly compact subring of $\widetilde{A}_{P}$, so that $L \supset \mathfrak{o}_{Q}$. $L$ spans $\widetilde{A}_{P} / K_{P}$ since it is open, and $L$ is finitely generated over $\mathfrak{o}_{P}$ by Lemma 1.1.

Conversely, if $E$ is a maximal order of $\tilde{A}_{P}$ for $\mathrm{o}_{P}$, then $E$ is trivially open and is linearly compact as a finitely generated space over $\mathrm{o}_{P}$. If $L \supset E$ is an open linearly compact subring, then by the previous argument $L$ is an order of $\tilde{A}_{P}$, since we used the maximality only to prove $L \supset \mathfrak{o}_{P}$. Therefore $L=E$, Q.E.D.

We now investigate the relation between the orders of $A / F(x)$ and those of $\widetilde{A}_{P} / F(x)_{Q}$. We state but do not prove

Lemma 2.1. If $J$ is an order of $A$ relative to $F[x]$ and $P$ a finite prime, then the closure $L$ of $J$ in $\tilde{A}_{P}$ is an open linearly compact subring of $\widetilde{A}_{P}$ containing $\mathfrak{D}_{Q}$.

Lemma 2.2. If $J$ is a maximal order of $A$ relative to $F[x]$ and $P$ a finite prime, then the closure $L$ of $J$ in $\widetilde{A}_{P}$ is a maximal order of $\tilde{A}_{P}$.

Proof. By the previous lemma, we know that $L$ is an order of $\tilde{A}_{P}$. To prove $L$ maximal, let $L^{\prime}$ be an order of $\widetilde{A}_{P}$ containing $L$. Following [5] and [12], we define $J^{\prime}$ as the set of all $a \in A \cap L^{\prime}$ such that for some integer $\alpha$, $p(x)^{\alpha} a \epsilon J$, where $p(x)$ is the irreducible polynomial from $F[x]$ giving rise to the prime divisor $Q$. Notice that $J^{\prime}$ is a subring of $A$ containing $J$ and hence spanning $A / F(x)$. Also, for some integer $\beta>0$, we have $J^{\prime} \subset p(x)^{-\beta} J$, since $J^{\prime}$ is contained in the linearly compact subspace $L^{\prime}$; thus $J^{\prime}$ is finitely generated over $F[x]$ as a submodule of $p(x)^{-\beta} J$. Therefore $J^{\prime}$ is an order of $A$; since $J$ is maximal, $J^{\prime}=J$.

Now suppose there is an element $a$ in $L^{\prime}$ not in $L$. Letting $u_{1}, \cdots, u_{m}$ be a basis of $A / F(x)$, we write $a=\sum a_{i} u_{i}, a_{i} \in F(x)_{Q}$. We can find elements $b_{i} \in F(x)$ close enough to the $a_{i}$ so that $b=\sum b_{i} u_{i}$ is also in $L^{\prime}$ but not in $L$, since $L$ is closed and $L^{\prime}$ is open. This $b$ is in $A$. There is a $g \in F[x]$ such that $g \cdot b \in J$, and we may factor $g$ as $g=p(x)^{\alpha} h$, where $h \in F(x)$ is prime to $p(x)$. Then $h b \in L^{\prime}$; and, since $h$ is a unit in $\mathfrak{o}_{Q}$ and $L$ is an $\mathfrak{o}_{Q}$-module,

[^3]$h^{-1} L \subset L$, which implies $h b \notin L . \quad$ But $h b$ is in $J^{\prime}=J \subset L$, a contradiction. Therefore $L^{\prime}=L$, Q.E.D.

Lemma 2.3. Let $u_{1}, \cdots, u_{m}$ be any basis of $A / F(x)$ and define, for each $P \in \mathfrak{M}, E_{P}=\mathfrak{o}_{Q} u_{1}+\cdots+\mathfrak{o}_{Q} u_{m}$. Then for almost all $P, E_{P}$ is a maximal order of $\widetilde{A}_{P}$ and equals $\mathfrak{o}_{P} u_{1}+\cdots+\mathfrak{o}_{P} u_{m}$.

Proof. Let $J$ be any maximal order of $A$ relative to $F[x], J_{P}$ the closure of $J$ in $\widetilde{A}_{P}$. By Lemma 2.2, $J_{P}$ is a maximal order of $\widetilde{A}_{P}$ for almost all $P$. Let $J=F[x] b_{1}+\cdots+F[x] b_{n}$ for some $b_{1}, \cdots, b_{n} \in A$, and let each $b_{i}=$ $\sum_{j} c_{i j} u_{j}$, where $c_{i j} \in F(x)$. Then the matrix ( $c_{i j}$ ) has no column consisting entirely of zeros. Since $J_{P}=\mathfrak{o}_{Q} b_{1}+\cdots+\mathfrak{p}_{Q} b_{n}$, it follows that $J_{P}=E_{P}$ for almost all $P$, Q.E.D.

Theorem 2.2. For each finite prime $P$, let $L_{P}$ be a maximal order of $\widetilde{A}_{P}$ such that almost all $L_{P}=\mathfrak{o}_{P} u_{1}+\cdots+\mathfrak{o}_{P} u_{m}$, where the $u$ 's form a basis of $A / F(x)$. Let $J$ denote the intersection of all $A \cap L_{P}$. Then $J$ is a maximal order of $A$ relative to $F[x]$.

Proof. Using the topological properties established up to now, one first proves that $J$ is an order of $A$. Then Lemma 2.2 implies that $J$ is a maximal order. We omit the details.

In order to clarify later parts of this paper, and to make a convenience rigorous, we now discuss isomorphisms between $\widetilde{A}$ and the matrix ring arising naturally from $\widetilde{A}$. That is, each $\widetilde{A}_{P}$ is a normal simple algebra over $K_{P}$ and is therefore algebraically isomorphic to the full matrix algebra $r_{P} \times r_{P} S_{P}$ over some skew-field $S_{P}$ with $K_{P}$ as center. If $A$ is isomorphic to the full matrix algebra $r \times r S_{0}$ over the skew-field $S_{0}$, then $K_{P} \otimes S_{0} \cong \mu_{P} \times \mu_{P} S_{P}$ for some integer $\mu_{P}$, so that $r_{P}=r \mu_{P}$. We define $\bar{S}_{P}$ to be $r_{P} \times r_{P} S_{P}$ and $\bar{S}$ to be the weak direct sum

$$
\bar{S}=\sum_{P \in \mathfrak{M}}^{\prime} \bar{S}_{P}
$$

in the sense that each matrix in $\bar{S}$ must have almost all its $P$-components taken from $r_{P} \times r_{P} \mathcal{O}_{P}$, in the notation of $\S 1 . \bar{S}$ can be given two topologies, one in which a fundamental system of neighborhoods of 0 consists of the subspaces of the form $X I$, where $X$ is a regular element of $\bar{S}$ and $I$ is the direct sum $\sum r_{P} \times r_{P} \mathcal{O}_{P}$, another which the direct-sum topology of each $\bar{S}_{P}$ gives rise to; that these topologies are actually the same is a consequence of Lemma 1.2.

Let us agree to denote $r_{P} \times r_{P} \mathfrak{O}_{P}$ by $I_{P}$ in what follows. An element $U$ of $\bar{S}$ (or of $\bar{S}_{P}$ ) is said to be unitary if it is a regular element of $\bar{S}$ (or of $\bar{S}_{P}$ ) such that both $U$ and $U^{-1}$ are in $I$ (or $I_{P}$ ). We shall need the decomposition [13, p. 107] of a regular matrix $C \in \bar{S}_{P}$ as

$$
\begin{equation*}
C=U\left(\delta_{i j} t^{\epsilon_{i}}\right) V \tag{2.2}
\end{equation*}
$$

where $U$ and $V$ are unitary in $\bar{S}_{P}, \delta_{i j}$ is the Kronecker delta, $t$ is a prime element for $P$, and the $e_{i}$ are rational integers.

We now have
Lemma 2.4. A maximal open linearly compact subring of $\bar{S}_{P}$ is always of the form $C^{-1} I_{P} C$, where $C$ is a regular element of $\bar{S}_{P}$.

For a proof, see Hasse, [5, pp. 519-520].
This lemma enables one to prove rather easily
Theorem 2.3. There is a topological $K$-isomorphism of $\tilde{A}$ onto $\bar{S}$. Any topological automorphism of $\bar{S}$ onto $\bar{S}$ transforms almost all $I_{P}$ onto themselves.

Now let $J_{1}$ and $J_{2}$ be maximal orders of $A$ relative to $F[x]$. Let $M$ be a left-ideal for $J_{1}$ and a right-ideal for $J_{2}$. Let $M_{P}$ denote the closure of $M$ in $\widetilde{A}_{P}$, and similarly for $J_{1}^{P}, J_{2}^{P}$. For each finite $P$ let $L_{P}$ be a maximal order of $\tilde{A}_{P}$ such that almost all $L_{P}$ are equal to $\mathfrak{o}_{P} u_{1}+\cdots+\mathfrak{o}_{P} u_{n}=E_{P}$, for a fixed basis $u_{1}, \cdots, u_{n}$ of $A / K$.

Theorem 2.4. For each finite $P, M_{P}$ is an open linearly compact left-module for $J_{1}^{P}$ and right-module for $J_{2}^{P}$; almost all $M_{P}=E_{P}$. Conversely, the intersection with $A$ of such $M_{P}$ 's is an ideal of $A$ relative to $F[x]$; in particular, our original ideal $M$ is the intersection of its components $M_{P}$. Each $M_{P}$ has the form $C_{P}^{-1} L_{P} C_{P}^{\prime}$, where $C_{P}$ and $C_{P}^{\prime}$ are regular elements of $\tilde{A}_{P}$ such that almost all $C_{P}^{-1} L_{P} C_{P}^{\prime}=L_{P}=E_{P}$.

We omit the proof of this theorem, as well as that of the following
Theorem 2.5. The maximal open linearly compact subrings of $\tilde{A}$ are all conjugate to each other in $\tilde{A}$. They are the direct sums of their P-components.

## 3. A general Riemann-Roch theorem

Let $V$ be a finite-dimensional vector space over the function field $K$. Let $V^{\prime}$ denote the dual space to $V$. As spaces of valuation vectors of $V$ and $V^{\prime}$, we define $\tilde{V}=R \otimes_{K} V$ and $\tilde{V}^{\prime}=R \otimes_{K} V^{\prime}$. The natural pairing $\left\langle v, v^{\prime}\right\rangle_{0}=$ $v^{\prime}(v)$ of $V$ and $V^{\prime}$ to $K$ can be extended uniquely by continuity to a pairing of $\tilde{V}$ and $\tilde{V}^{\prime}$ to $R$ : $\left\langle\tilde{v}, \tilde{v}^{\prime}\right\rangle$ is a continuous, $R$-bilinear map of $\tilde{V} \times \tilde{V}^{\prime}$ into $R$. Letting $\chi$ be an admissible character of $R$, we define

$$
\begin{equation*}
\left[v, v^{\prime}\right]=\chi\left(\left\langle v, v^{\prime}\right\rangle\right), \quad v \in \tilde{V}, \quad v^{\prime} \in \tilde{V}^{\prime} \tag{3.1}
\end{equation*}
$$

and obtain thereby a continuous $F$-bilinear map of $\tilde{V} \times \tilde{V}^{\prime}$ into $F$.
Theorem 3.1. The mapping $f: \tilde{V} \rightarrow X\left(\tilde{V}^{\prime}\right)$ of $\tilde{V}$ to the character space of $\tilde{V}^{\prime}$ given by $f(\tilde{v})=[\tilde{v}, \quad]$ is a topological isomorphism onto; that is, the dual pairing (3.1) is topological. Under this pairing, the annihilator $A(V)$ of $V$ is $V^{\prime}$.

Proof. For any basis $\left\{v_{i}\right\}$ of $V / K$, let $\left\{v_{i}^{\prime}\right\}$ be the dual basis of $V^{\prime} / K$. With respect to these (or any) bases, $\tilde{V}$ and $\tilde{V}^{\prime}$ are immediately seen to be paired as direct sums of copies of the self-dual space $R$ (see §1). But [ , ] becomes this very pairing when expressed in terms of these bases. This ob-
servation also shows immediately that $A(V)=V^{\prime}$, since $A(K)=K$ under the pairing of $R$ to itself mentioned in §1, Q.E.D.

If $M$ is an open linearly compact subspace of $\tilde{V}$, then $M \cap V$ is finitedimensional over $F$, by Lemma 1.1. Thus we may define for such $M$

$$
l(M)=\operatorname{dim}_{F}(M \cap V)
$$

and similarly for such subspaces of $\tilde{V}^{\prime}$.
We introduce the unique function $\nu(, \quad)$ defined on ordered pairs of open linearly compact subspaces of $\tilde{V}$ such that

$$
\begin{array}{ll}
\text { (i) } & \nu\left(M_{1}, M_{2}\right)+\nu\left(M_{2}, M_{3}\right)=\nu\left(M_{1}, M_{3}\right)  \tag{3.2}\\
\text { (ii) } & \nu\left(M_{1}, M_{2}\right)=\operatorname{dim}_{F}\left(M_{1} / M_{2}\right) \quad \text { if } \quad M_{1} \supset M_{2} .
\end{array}
$$

The existence and uniqueness of this $\nu$-function are quite easy to prove; ${ }^{9}$ and it follows that

$$
\begin{equation*}
\nu\left(M_{1}, M_{2}\right)=\nu_{1}\left(\pi M_{1}, \pi M_{2}\right)+\nu_{2}\left(M_{1} \cap V, M_{2} \cap V\right) \tag{3.3}
\end{equation*}
$$

where $\nu_{1}$ and $\nu_{2}$ are the analogous $\nu$-functions for the spaces $\widetilde{V} / V$ and $V$, respectively, and $\pi$ is the natural map from $\tilde{V}$ onto $\tilde{V} / V$.

Now let $M^{\prime}$ be the annihilator with respect to the dual pairing (3.1) of the divisor $M$ of $\tilde{V} . \quad M^{\prime}$ is open in $\tilde{V}^{\prime}$ by the continuity of [ , ] and is linearly compact as the dual space to the discrete space $\tilde{V} / M$. By Theorem 3.1, the annihilator of $M \cap V$ is the (closed) subspace $M^{\prime}+V^{\prime}$ of $\tilde{V}^{\prime}$; therefore

$$
\begin{equation*}
l(M)=\operatorname{dim}_{F}\left(\tilde{V}^{\prime} /\left(M^{\prime}+V^{\prime}\right)\right) \tag{3.4}
\end{equation*}
$$

since $\tilde{V}^{\prime} /\left(M^{\prime}+V^{\prime}\right)$ is dual to the finite-dimensional space $M \cap V$.
We shall now restate (3.3) in terms of $l$. We have

$$
\begin{aligned}
\nu_{1}\left(\pi M_{1}, \pi M_{2}\right) & =-\nu_{1}\left(\tilde{V} / V, \pi M_{1}\right)+\nu_{1}\left(\tilde{V} / V, \pi M_{2}\right) \\
& =-\operatorname{dim}_{F}\left(\tilde{V} /\left(M_{1}+V\right)\right)+\operatorname{dim}_{F}\left(\tilde{V} /\left(M_{2}+V\right)\right) \\
& =-l\left(M_{1}^{\prime}\right)+l\left(M_{2}^{\prime}\right)
\end{aligned}
$$

by (3.2) and (3.4). Also, $\nu_{2}(M \cap V, 0)=l(M)$, so that we may put (3.3) as

Theorem 3.2. For any two open linearly compact subspaces $M_{1}, M_{2}$ of $\tilde{V}$, we have

$$
\nu\left(M_{1}, M_{2}\right)=l\left(M_{1}\right)-l\left(M_{1}^{\prime}\right)-\left(l\left(M_{2}\right)-l\left(M_{2}^{\prime}\right)\right)
$$

where $M_{1}^{\prime}$ and $M_{2}^{\prime}$ are the annihilators of $M_{1}$ and $M_{2}$, respectively, with respect to the dual pairing (3.1).

Let $A$ and $B$ be simple algebras of finite rank over the center $K$. Assume now that $V$ is a unitary left $A$-, right $B$-module. Then $V^{\prime}$ is naturally a left $B$-, right $A$-module, and the pairing $\langle\quad, \quad\rangle_{0}$ satisfies $\left\langle a v b, v^{\prime}\right\rangle_{0}=\left\langle v, b v^{\prime} a\right\rangle_{0}$ for all $a \in A, b \in B, v \in V$, and $v^{\prime} \in V^{\prime}$.

[^4]There is a unique way to make $\tilde{V}$ a unitary left $\tilde{A}-$, right $\widetilde{B}$-module (with the action denoted for the moment by a dot) such that

$$
\begin{array}{cr}
s \cdot v=v \cdot s=s \otimes v, & s \in R, \quad v \in V, \\
a \cdot v=a v, \quad v \cdot b=v b, & a \in A, \quad b \in B, \quad v \in V .
\end{array}
$$

The map of $\tilde{A} \times \tilde{V} \times \tilde{B}$ to $\tilde{V}$ which sends ( $\tilde{a}, \tilde{v}, \tilde{b}$ ) into $\tilde{a} \cdot \tilde{v} \cdot \tilde{b}$ is then continuous. The analogous result holds for $\tilde{V}^{\prime}$. Furthermore, there is a unique $R$-bilinear pairing $\langle$,$\rangle of \widetilde{V}$ and $\widetilde{V}^{\prime}$ to $R$ such that $\left\langle v, v^{\prime}\right\rangle=\left\langle v, v^{\prime}\right\rangle_{0}$ for $v \in V, v^{\prime} \in V^{\prime} .\langle, \quad\rangle$ is then a continuous map of $\tilde{V} \times \tilde{V}^{\prime}$ into $R$ such that $\left\langle\tilde{a} \tilde{v} \tilde{b}, \tilde{v}^{\prime}\right\rangle=\left\langle\tilde{v}, \tilde{b} \tilde{v}^{\prime} \tilde{a}\right\rangle$ for $\tilde{a} \in \tilde{A}, \tilde{v} \in \tilde{V}$, etc. Finally, we have

Lemma 3.1. If $A$ acts faithfully on $V$, then so does $\widetilde{A}_{P}$ on $\widetilde{V}_{P}$.
Proof. The result follows immediately from the general structure theorem for such modules (see [2, p. 46]), which says that if $A$ is $K$-isomorphic to the full matrix ring $r \times r S_{0}$ over the skew-field $S_{0}$, then $V$ is $K$-isomorphic to $r \times r_{1} S_{0}$, for some $r_{1}$, and the action of $A$ on $V$ is given by the usual matrix multiplication. On tensoring with $K_{P}$ we get the desired result.

We shall now prove a formula which, in some cases, allows us to compute $\nu\left(M_{1}, M_{2}\right)$. Let $M_{2}$ be a given open linearly compact subspace of $\widetilde{V}$; suppose furthermore that $M_{2}$ is an $\mathfrak{0}$-module, which implies in particular that it is the direct sum of its $P$-components. Then the same properties hold for the subspace $M_{1}=a^{-1} M_{2} b$, where $a$ and $b$ are regular elements of $\widetilde{A}$ and $\widetilde{B}$, respectively. We shall compute $\nu\left(M_{1}, M_{2}\right)$.

First assume $M_{1} \supset M_{2}$. Then $\nu\left(M_{1}, M_{2}\right)=\operatorname{dim}_{F} M_{1} / M_{2}$, and it will suffice to compute the $F$-dimension $\rho_{P}$ of the $P$-component $M_{1 P} / M_{2 P}$; for then $\nu\left(M_{1}, M_{2}\right)=\sum_{P} \rho_{P}$.

Lemma 3.2. Let the subspace $M_{P}$ of $\tilde{V}_{P}$ be an open, linearly compact $\mathfrak{o}_{P^{-}}$ module. Then there exists a basis $\left\{w_{j}\right\}$ of $\widetilde{V}_{P} / K_{P}$ such that $M_{P}=\sum \mathfrak{o}_{P} w_{j}$.

Proof. The lemma follows immediately from part 1 of [13, §108] when we observe that the openness of $M_{P}$ implies that $M_{P}$ contains a basis of $\widetilde{V}_{P} / K_{P}$.

We can now proceed to compute $\rho_{P}$, on the assumption that $M_{1 P}=$ $\left(a^{-1} M_{2} b\right)_{P} \supset M_{2 P}$. Let $w_{1}, \cdots, w_{n}$ be the basis of $\widetilde{V}_{P} / K_{P}$ contained in $M_{2 P}$ described in Lemma 3.2. With respect to this basis, the operation at $P$ of $a^{-1}$ and $b$ leads to a nonsingular $n \times n$ matrix $\left(\alpha_{i j}\right)$ over $K_{P}$ as follows:

$$
a_{P}^{-1} w_{i} b_{P}=\sum_{j=1}^{n} \alpha_{i j} w_{j}, \quad a_{i j} \in K_{P}, \quad i=1, \cdots, n
$$

Then

$$
M_{1 P}=\sum_{i} \mathfrak{o}_{P}\left(a_{P}^{-1} w_{i} b_{P}\right)=\left\{\sum_{i, j} \mathfrak{a}_{i} \alpha_{i j} w_{j} ; \mathfrak{a}_{i} \in \mathfrak{o}_{P}\right\}
$$

If we express $M_{1 P}$ as a set of $n$-tuples, the coefficients with respect to the basis $w_{1}, \cdots, w_{n}$, we find

$$
M_{1 P}=\left(\mathfrak{o}_{P}, \cdots, \mathfrak{o}_{P}\right)\left(\alpha_{i j}\right)=\left(1 \times n \mathfrak{o}_{P}\right)\left(\alpha_{i j}\right)
$$

In the same way $M_{2 P}$ becomes simply $1 \times n 0_{P}$. Thus

$$
\rho_{P}=\operatorname{dim}_{F}\left(1 \times n \mathfrak{o}_{P}\right)\left(\alpha_{i j}\right) /\left(1 \times n \mathfrak{o}_{P}\right)
$$

In order to calculate $\rho_{P}$, we represent $\left(\alpha_{i j}\right)$ according to (2.2) as $\left(\alpha_{i j}\right)=$ $u\left(\delta_{i j} \tau^{e_{i}}\right) v$, where $u$ and $v$ are unitary $n \times n$ matrices over $K_{P}$ and $\nu_{P}(\tau)=1$. Now since $\left(1 \times n \mathbf{o}_{P}\right) u=\left(1 \times n \mathfrak{o}_{P}\right)$, our above factor-space is $\left(1 \times n \mathfrak{o}_{P}\right)\left(\delta_{i j} \tau^{e_{i}}\right) v /\left(1 \times n \mathfrak{o}_{P}\right)$, which is isomorphic over $F$ to $\left(1 \times n \mathfrak{o}_{P}\right)$.

$$
\left(\delta_{i j} \tau^{e_{i}}\right) /\left(1 \times n \mathfrak{o}_{P}\right)=\left(\mathfrak{p}^{e_{1}}, \cdots, \mathfrak{p}^{e_{n}}\right) /\left(\mathfrak{o}_{P}, \cdots, \mathfrak{o}_{P}\right) \cong \Omega_{P} \oplus \cdots \oplus \Omega_{P}
$$

with $-\left(e_{1}+\cdots+e_{n}\right)$ summands (the $e_{i}$ being all nonpositive because of the assumption that $M_{1} \supset M_{2}$ ). Therefore the desired dimension $\rho_{P}$ is $-n_{P}\left(e_{1}+\cdots+e_{n}\right)$. And this in turn may be written

$$
\rho_{P}=n_{P} \nu_{P}\left(\operatorname{det}\left(\alpha_{i j}\right)^{-1}\right)
$$

To find the relation between $\operatorname{det}\left(\alpha_{i j}\right)$ and the norms of $a$ and $b$, we use the structure theorem [2, p. 46] quoted in the proof of Lemma 3.1. Here it is easy to see that the matrix $\left(\beta_{i j}\right)$ for $a \in A$ arising out of $a v_{i}=\sum \beta_{i j} v_{j}$, with $\beta_{i j} \epsilon K,\left\{v_{j}\right\}$ a basis of $V / K$, is the $r_{1}$-fold repetition of the matrix obtained when one replaces each $S_{0}$-entry of $a$ with its matrix in the regular representation of $S_{0} / K$. (The positive integer $r_{1}$ is the number of columns in the isomorph of $V$ given by the structure theorem.) Since the matrix of $a$ in the regular representation of $A$ over $K$ is the $r$-fold repetition of the same matrix, it follows immediately that $\operatorname{det}\left(\beta_{i j}\right)=N(a)^{r_{1} / r}$. When we pass to the local situation, both $r_{1}$ and $r$ are multiplied by $\mu_{P}$, leaving the exponent unchanged.

Similarly, the analogous matrix $\left(\gamma_{i j}\right)$ for $b \in B$ satisfies $\operatorname{det}\left(\gamma_{i j}\right)=N^{\prime}(b)^{r_{2} / r^{\prime}}$ for some positive integer $r_{2}$ (the number of rows in the appropriate isomorph of $V$, when $B$ is isomorphic to $r^{\prime} \times r^{\prime} S_{0}^{\prime}, S_{0}^{\prime}$ a normal skew-field over $K$ ). $N^{\prime}$ denotes the norm of the regular representation of $B / K$. It can be extended uniquely to $\widetilde{B} / R$.

Returning now to the provocation for all this, we see that there exist rational numbers $\rho_{1}=r_{1} / r$ and $\rho_{2}=r_{2} / r^{\prime}$ depending only on the module structure of $(A, V, B)$ such that

$$
\operatorname{det}\left(\alpha_{i j}\right)^{-1}=N_{P}\left(a_{P}\right)^{\rho_{1}} N_{P}\left(b_{P}\right)^{-\rho_{2}}
$$

Therefore $\rho_{P}=n_{P}\left(\rho_{1} \nu_{P} N_{P}\left(a_{P}\right)-\rho_{2} \nu_{P} N_{P}\left(b_{P}\right)\right)$. And now we can assert
Lemma 3.3. If $M_{0}$ is an open, linearly compact $\mathfrak{o}-m o d u l e ~ c o n t a i n e d ~ i n ~ \tilde{V}$, and if $M=a^{-1} M_{0} b$ contains $M_{0}$, for regular $a \in \widetilde{A}, b \in \widetilde{B}$, then there exist positive rational numbers $\rho_{1}$ and $\rho_{2}$ depending only on the module structure of ( $A, V, B$ ) such that

$$
\nu\left(M, M_{0}\right)=\operatorname{dim}_{F}\left(M / M_{0}\right)=n\left(N(a)^{\rho_{1}} N^{\prime}(b)^{-\rho_{2}}\right)
$$

where $n(c)$ denotes the degree of the regular element $c$ in $R$.
Proof. We need only recall that $\operatorname{dim}_{F} M / M_{0}=\sum \rho_{P}$ and that the degree of a regular element $c$ of $R$ is $\sum n_{P} \nu_{P}\left(c_{P}\right)$.

Now we need to consider how to compute $\nu\left(M, M_{0}\right)$ in the general case, when $M$ may not contain $M_{0}$. We still assume that $M=a^{-1} M_{0} b$ for regular $a \in \widetilde{A}, b \in \widetilde{B}$, however, and that $M_{0}$ is an open, linearly compact $\mathfrak{o}$-module. We simply define $\nu\left(M, M_{0}\right)=n\left(N(a)^{\rho_{1}} N^{\prime}(b)^{-\rho_{2}}\right)$; the verification that property (i) of (3.2) holds for this $\nu$-function is routine, and we have just proved that (ii) holds. Therefore we have proved

Theorem 3.3. Let $A$ and $B$ act faithfully on $V$, and let the subspace $M_{0}$ of $\tilde{V}$ be an open, linearly compact $\mathrm{D}-\mathrm{module}$. Then for all subspaces of the form $M=a^{-1} M_{0} b$, with a regular in $\widetilde{A}, b$ regular in $\widetilde{B}$, the $\nu$-function (3.2) satisfies

$$
\nu\left(M, M_{0}\right)=n\left(N(a)^{\rho_{1}} N(b)^{-\rho_{2}}\right)
$$

for certain positive rational numbers $\rho_{1}$ and $\rho_{2}$ depending only on the structure of the module $(A, V, B)$ and not on $M_{0}$ or $M$.

Under the notations of Theorem 3.3, let us define the degree of $M$ with respect to $M_{0}$ as

$$
\begin{equation*}
n(M)=n\left(N(a)^{\rho_{1}} N^{\prime}(b)^{-\rho_{2}}\right) \tag{3.5}
\end{equation*}
$$

These subspaces $M$ will play the role of divisors in our generalization of the Riemann-Roch theorem to the module $V$.

Let the annihilator of $\mathfrak{o}$ in $R$ with respect to $\chi$ be $d^{-1} \mathfrak{o}$ for a regular $d \epsilon R$. For a given $P \in \mathfrak{M}$, let $M_{0 P}=\mathfrak{o}_{P} w_{1}+\cdots+\mathfrak{o}_{P} w_{n}$ in accordance with Lemma 3.2. Then the $P$-component of the annihilator of $M_{0}$ is $M_{0 P}^{\prime}=d_{P}^{-1} \mathfrak{0}_{P} w_{1}^{\prime}+$ $\cdots+d_{P}^{-1} \mathrm{o}_{P} w_{n}^{\prime}$, where $w_{1}^{\prime}, \cdots, w_{n}^{\prime}$ is the dual basis to $w_{1}, \cdots, w_{n}$. And that for $M$ is $M_{P}^{\prime}=d_{P}^{-1} \mathfrak{o}_{P}\left(b_{P}^{-1} w_{1}^{\prime} a_{P}\right)+\cdots+d_{P}^{-1} \mathfrak{o}_{P}\left(b_{P}^{-1} w_{n}^{\prime} a_{P}\right)$. Thus Lemma 3.2 and Theorem 3.3 allow us to state Theorem 3.2 in a more explicit form, which we call our generalization of the Riemann-Roch Theorem to $V$ :

Theorem 3.4. Under the hypotheses and notations of Theorem 3.3, we have

$$
l(M)=l\left(M^{\prime}\right)+n(M)-\left(l\left(M_{0}^{\prime}\right)-l\left(M_{0}\right)\right)
$$

where $n(M)$ is the degree of $M$ with respect to $M_{0}$ as defined in (3.5). Here the quantity $l\left(M_{0}^{\prime}\right)-l\left(M_{0}\right)$ depends only on $M_{0}$, not on the module structure.

Remark. The functions $l$ and $n$ appearing here are "class-functions"; that is, if $a$ and $b$ are regular elements of $A$ and $B$, respectively, then $l(a M b)=$ $l(M)$ and $n(a M b)=n(M)$, for any open, linearly compact $\mathfrak{o}$-module $M$.

This theorem contains the classical Riemann-Roch theorem: We take $A=V=B=K, M_{0}=\mathfrak{o}, M=a^{-1} \mathfrak{v}$ for a regular element $a$ of $R$; then $l(\mathfrak{D})=1, l\left(\mathfrak{o}^{\prime}\right)=g$, the genus of $K$, and the degree of $M=a^{-1} \mathfrak{o}$ with respect to $M_{0}=\mathfrak{o}$ is $n\left(a^{-1} \mathfrak{v}\right)=n(a)$ in our definition. For the classical theorem we define $l^{*}(a)=l\left(a^{-1} \mathrm{D}\right)$ and obtain from our theorem the classical form of the Riemann-Roch theorem, namely,

$$
l^{*}(a)=l^{*}\left(a^{\prime}\right)+n(a)-g+1
$$

where by $a^{\prime}$ we understand any regular $b \in R$ such that $\left(a^{-1} \mathfrak{D}\right)^{\prime}=b 0$; it is well known that $b=a^{-1} d$, when $\mathfrak{v}^{\prime}=d^{-1} \mathrm{o}$.

In the classical case, divisors are defined as the set of all $a \mathrm{D}, a$ regular in $R$; they can be characterized as the set of all open, linearly compact $\mathfrak{o}$-modules. For the module $V$ considered in this paper we have taken for divisors the set of all subspaces of $\widetilde{V}$ which are open, linearly compact $\mathfrak{D}$-modules. In the general Theorem 3.2 the algebras $A$ and $B$ play no role, nor is the assumption that the subspaces be 0 -modules needed there. The use of the algebras $A$ and $B$ and the assumption of closure under $\mathbf{o}$-multiplication is that they enable us to find a nice formula for the $\nu$-function, provided one of the two divisors can be obtained from the other via multiplication by regular elements of $\widetilde{A}$ and $\widetilde{B}$. As noted above, this relation holds between any two divisors in the classical case; it does not hold in general for our module $V$, however. But in one example where this relation fails to be universal, $V$ is not irreducible as a double module; in another such example, $B$ is not normal as a simple algebra over $K$; but in these examples divisors are further restricted to be modules with respect to the actions of maximal open linearly compact subrings of $\widetilde{A}$ and $\widetilde{B}$. Whether the reasonable assumptions of faithfulness, irreducibility as a double module, and normality imply that any (reasonably defined) divisor can be obtained from any other divisor via multiplication by regular elements of $\widetilde{A}$ and $\widetilde{B}$ is an open question. This relation between pairs of divisors is an equivalence relation, and Theorem 3.4 holds for divisors taken from any one class.

We now turn to a situation where, when the notion of divisor is suitably restricted, the relation in question holds between any two divisors.

## 4. Simple algebras over $K$

Let $A=V=B$ be a simple algebra over the center $K$, with the action being multiplication in $A$. Letting $T$ denote the reduced trace from $A$ to $K$, we pair $A$ to itself by setting $\left\langle c, c^{\prime}\right\rangle_{0}=T\left(c c^{\prime}\right)$, for $c, c^{\prime} \in A$. Then $\left\langle a c b, c^{\prime}\right\rangle_{0}=$ $\left\langle c, b c^{\prime} a\right\rangle_{0}$ for all $a, b \in A$. Our dual pairing of $\widetilde{A}$ to itself becomes

$$
\begin{equation*}
\left[c, c^{\prime}\right]=\chi T\left(c c^{\prime}\right), \quad c, c^{\prime} \in \widetilde{A} \tag{4.1}
\end{equation*}
$$

In order to determine the numbers $\rho_{1}$ and $\rho_{2}$ appearing in the formula (3.5), we need only recall that if $A$ is isomorphic to $r \times r S_{0}, S_{0}$ being a skew-field with center $K$, then $\rho_{1}=r_{1} / r$, where $r_{1}$ is the number of columns of $r \times r_{1} S_{0}$, the isomorph of $V$. Thus $r_{1}=r$ and $\rho_{1}=1$. Similarly $\rho_{2}=1$. Therefore, if $M_{0}$ is a divisor of $\tilde{A}$ and if $M=a^{-1} M_{0} b$, for regular elements $a, b \in \tilde{A}$, then

$$
\begin{equation*}
n(M)=\nu\left(M, M_{0}\right)=n\left(N\left(a b^{-1}\right)\right) \tag{4.2}
\end{equation*}
$$

where $N$ denotes the norm (of the regular representation) from $\tilde{A}$ to $R$. ( $N$ is the "direct product" of the local norms from $\tilde{A}_{P}$ to $K_{P}$.)

Divisors. For convenience we shall denote the set of all maximal open linearly compact subrings of $\widetilde{A}$ as $\mathcal{L}(\widetilde{A})=\mathscr{L}$. In the present situation we
restrict our divisors to be open, linearly compact subspaces $M$ of $A$ such that there exist $H, J \in \mathcal{L}$ for which $H M J=M$. For this divisor we shall occasionally denote $H$ by $M_{l}$ and $J$ by $M_{r}$.

We can easily prove that all divisors are equivalent in the sense discussed at the end of $\S 3$ by using the isomorphism of $\S 2$ between $\widetilde{A}$ and the appropriate space of matrices $\bar{S}$ : The members of $\mathfrak{L}(\bar{S})$ are of the form $C^{-1} I C$, where $C$ is a regular element of $\bar{S}$ and $I$ is the subring of "integral" matrices defined near the end of $\S 2$. If $C_{1}^{-1} I C_{1}$ and $C_{2}^{-1} I C_{2}$ are two such subrings equalling $M_{1 l}$ and $M_{1 r}$ for a divisor $M_{1}$, then $C_{1} M_{1} C_{2}^{-1}=M_{2}$ is a two-sided $I$-module; it is easy to prove that $M_{2}$ must then be of the form $c I$, where $c$ is a regular element of $\bar{S}$ of the form

$$
\begin{equation*}
c_{P}=\left(\delta_{i j} a_{P}\right), \quad a_{P} \in S_{P} \tag{4.3}
\end{equation*}
$$

for each $P \epsilon \mathfrak{M}$. Thus $M_{1}=C_{1}^{-1} c I C_{2}$. Conversely, if $C_{1}$ and $C_{2}$ are regular elements of $\bar{S}$, then $C_{1}^{-1} I C_{2}$ is a divisor of $\bar{S}$. If we transfer these results back to $A$ by means of our isomorphism, we can now assert

Lemma 4.1. Each divisor of $\tilde{A}$ is of the form $a^{-1} H b$ for regular $a, b \in \tilde{A}$ and for a fixed subring $H \in \mathcal{L}$. Conversely, every such subspace of $\widetilde{A}$ is a divisor of $\widetilde{A}$.

For later use, we shall now discuss the element $c$ mentioned above. First we state a criterion for equality between divisors, the proof of which follows rapidly from the decomposition (2.2):

Lemma 4.2. The divisors $C^{-1} I C_{0}$ and $C^{\prime-1} I C_{0}^{\prime}$ of $\bar{S}$ are equal if and only if there is an element $a \in \bar{S}$ locally of the form (4.2) such that both $a C^{\prime} C^{-1}$ and $a C_{0}^{\prime} C_{0}^{-1}$ are unitary.

In our proof of Lemma 4.1 we saw that any divisor $M_{2}$ of $\bar{S}$ which is a two-sided $I$-module has the form $M_{2}=c I$, where each $P$-component of $c$ is in $S_{P}$, i.e., is a diagonal matrix. The above lemma implies that $c$ is uniquely determined up to a unitary factor. Furthermore, $c I=I c$; and $c$ belongs to $\bar{S}_{0}$, the image of $\bar{S}_{0}$ in $\bar{S}$. If $H \in \mathcal{L}$, let us define an $H$-unit as a regular element $b$ of $\tilde{A}$ such that $b \in H$ and $b^{-1} \in H$.

Then we have proved
Lemma 4.3. If the divisor $M$ of $\tilde{A}$ is a two-sided $H$-module, for $H \in \mathbb{L}$, then $M$ is of the form $c H=H c$, where $c$ is a regular element of $\widetilde{S}_{0}$ uniquely determined up to an H-unit factor.

This lemma assumes that one selects a particular $S_{0} \subset A$.
The degree of a given divisor $M$ satisfies, when $H \in \mathscr{L}$,

$$
\nu(M, H)=\nu\left(M, a^{-1} H a\right)
$$

for all regular $a \epsilon \tilde{A}$; or, in other words, since all $H \in \mathfrak{L}$ are conjugate to each other, the degree of a divisor is invariant with respect to the members of this class 2. From now on in this section we restrict our definition of the degree
(3.5) of a divisor by requiring that $M_{0}$ belong to $\mathfrak{\&}$. (It is true that the degree of a divisor is invariant with respect to the members of any conjugacyclass of divisors, but $\mathfrak{L}$ is a naturally distinguished such class.)

Let $H_{0}$ and $H$ be any two members of $\mathfrak{£}$. As divisors, these have degree 0 , so that if we apply Theorem 3.2 to $H$ and $H_{0}$, we find that $l\left(H^{\prime}\right)-l(H)=$ $l\left(H_{0}^{\prime}\right)-l\left(H_{0}\right)$, which proves that this quantity is an invariant of $A$. We denote it by

$$
\begin{equation*}
\xi(A)=l\left(H^{\prime}\right)-l(H), \quad H \in \mathscr{L} \tag{4.4}
\end{equation*}
$$

If $A$ is isomorphic to $r \times r S_{0}$, as before, then we can prove
Proposition 4.1. $\xi(A)=r^{2} \xi\left(S_{0}\right)$.
Proof. We use the isomorphism of $\widetilde{A}$ with $\bar{S}$ : Since $I \in \mathscr{L}(\bar{S})$, it suffices to consider $l\left(I^{\prime}\right)-l(I)$; now $I=r \times r \mathcal{O}$, and $\mathcal{O} \epsilon \mathscr{L}\left(\bar{S}_{0}\right) ; I^{\prime}=r \times r \mathcal{O}^{\prime}$; therefore $\xi(A)=r^{2} \xi\left(S_{0}\right)=r^{2}\left(l\left(\mathcal{O}^{\prime}\right)-l(\mathcal{O})\right)$. We shall discuss some questions related to this invariant farther on.

Canonical divisors. In the classical case $A=K$, the canonical divisor corresponding to a given admissible character $\chi$ is defined as the inverse of the annihilator of $\mathfrak{o}$ with respect to the dual pairing $[c, b]=\chi(c b), c, b \in R$. That is, the annihilator $\mathfrak{v}^{\prime}$ equals $d^{-1} \mathrm{v}$ for some regular $d \epsilon R$; and $d \mathrm{o}$ is the corresponding canonical divisor. As $a$ runs through the nonzero elements of $K$, all admissible characters of $R$ are obtained as $\chi_{a}$, where $\chi_{a}(\quad)=\chi(a(\quad))$; the corresponding canonical divisor is $a d o$. Here the divisors form an abelian group and the principal divisors (namely, $a \mathfrak{D}, 0 \neq a \in K$ ) a subgroup $\Pi$. Thus, modulo $\Pi$, the set of all canonical divisors is precisely a coset. Analogous conditions hold in the general case.

Let $M=a^{-1} H b$ be a divisor of $\tilde{A}$. By the inverse of $M$ we simply mean the divisor $M^{-1}=b^{-1} H a$. The relations $M M^{-1}=M_{l}$ and $M^{-1} M=M_{r}$ hold.

In the present case, "the" ring of integers in $\tilde{A}$ is determined only up to inner automorphism. Each $H \in \mathscr{L}$ will have a collection of canonical divisors; those of $a^{-1} H a$ will be the conjugates under $a$ of those of $H$.

Let $H \in \mathscr{L}$, and consider the annihilator $H^{\prime}$ of $H$ with respect to (4.1). $H^{\prime}$ is an open, linearly compact, two-sided $H$-module; in other words $H^{\prime}$ is a divisor of $\widetilde{A}$. We call $H^{\prime-1}$ the canonical divisor for $H$ corresponding to $\chi$. By Lemma 4.3, $H^{\prime}$ is of the form $\mathfrak{D}^{-1} H$, for some regular $\mathfrak{D} \epsilon S_{0}$ such that $\mathfrak{b} H=H \mathfrak{b}$; and $\mathfrak{b}$ is uniquely determined up to an $H$-unit factor. If $a$ is a regular element of $\tilde{A}$, then the canonical divisor for $a^{-1} H a$ with respect to $\chi$ is $a^{-1} \searrow H a$, or $a^{-1} \searrow a \cdot a^{-1} H a$. Therefore we need only determine the canonical divisors for $H$.

By Theorem 3.1, the annihilator of $A$ is $A$ itself; therefore we obtain all characters of $\tilde{A}$ which vanish on $A$ as $\chi_{a}=[a],, a \in A$. Among these, the ones which give rise to dual pairings of $\widetilde{A}$ with itself are those with $a$ regular. For such $a$, the pairing is $\left[c_{1}, c_{2}\right]_{a}=\chi T\left(a c_{1} c_{2}\right)$. When $a \notin K$, this pairing is not symmetric, so we must speak of left- and right-hand annihilators and
canonical divisors. For example, the left-hand annihilator $A_{l}(H)$ of $H$ is the set of all $c_{1}$ such that $\left[c_{1}, H\right]_{a}=0$; or $A_{l}(H)=a^{-1} H^{\prime}=a^{-1} b^{-1} H$. (We have now $A_{r}\left(A_{l}(H)\right)=H$, but $A_{l}\left(A_{l}(H)\right)=a^{-1} H a$.) The corresponding canonical divisor is $\grave{\searrow} H a$, which we shall call the right-hand canonical divisor.

We now define, by analogy with the case $A=K$, equivalence relations $\sim_{l}, \sim_{r}$ on the set of all divisors of $\widetilde{A}$. If $M_{1}$ and $M_{2}$ are divisors of $\widetilde{A}$, we say
$M_{1} \sim_{l} M_{2}\left(M_{1} \sim_{r} M_{2}\right)$ if and only if there exists a regular $a \in A$ such that $M_{1}=a M_{2}\left(M_{1}=M_{2} a\right)$.

We see that the set of left-hand canonical divisors for $H$ is precisely the equivalence-class containing $\delta H$ under $\sim_{l}$, and correspondingly for the righthand ones. If $a$ and $c$ are regular elements of $A$ and $\widetilde{A}$, respectively, then the right-hand canonical divisor for $c^{-1} H c$ corresponding to [ , $]_{a}$, is $c^{-1} \delta H c a$.

If $M=a^{-1} H b, H \in \mathscr{L}$, is an arbitrary divisor of $\tilde{A}$, then $M^{\prime}=b^{-1} H^{\prime} a=$ $b^{-1} \mathfrak{D}^{-1} H a$. If we (inadequately) denote $b^{-1} \mathfrak{b} b$ by $\mathfrak{D}_{M}$, then $M^{\prime}=\mathfrak{D}_{M}^{-1} M^{-1}$. We can also say that $M^{\prime}=\left(b^{-1} \delta^{-1} H b\right) M^{-1}$, the first factor being the inverse of the canonical divisor for $b^{-1} \mathrm{Hb}$. The left-annihilator of $M$ with respect to [ , $]_{c}, c$ regular in $A$, is $c^{-1} b^{-1} b^{-1} H a$.

We can now assert our generalization of the Riemann-Roch theorem for $A$ :
Theorem 4.1. If $M$ is any divisor of $\tilde{A}$, then

$$
l(M)=l\left(\mathfrak{D}_{M}^{-1} M^{-1}\right)+n(M)-r^{2} \xi\left(S_{0}\right) .
$$

Corollary 4.1. All canonical divisors have degree

$$
-2 r^{2} \xi\left(S_{0}\right)=n\left(N(\mathfrak{d})^{-1}\right)=-n\left(H^{\prime}\right), \quad H^{\prime} \in \mathscr{L}
$$

Proof. That they all have the same degree $n\left(N(\mathfrak{d})^{-1}\right)$ follows from the construction of them in the preceding paragraphs and the fact that $K$-elements have degree 0 in $R$. We obtain the corollary by putting $M=H^{\prime}, H \in \mathscr{L}$, in the theorem.

We can also derive as a corollary the generalization of Riemann's theorem in the special case $A=S_{0}$.

Corollary 4.2. If $M$ is a divisor of $\widetilde{S}_{0}$ with degree $n(M)>2 \xi\left(S_{0}\right)$, then

$$
l(M)=n(M)-\xi\left(S_{0}\right)
$$

Proof. We first notice that for any divisor $M, c \in M \cap A$ implies $n(N c) \geqq$ $-n(M)$. Therefore, $n(M)<0$ implies $l(M)=0$. Now if our given $M$ is of the form $a^{-1} H b$, then we see, from $M^{\prime}=b^{-1} b^{-1} H a$, that $n\left(M^{\prime}\right)=$ $n(N D)-n(M)<0$. Therefore, $l\left(M^{\prime}\right)=0$.

Remark 4.1. Theorem 4.1 is usually stated in terms of the functions $l^{*}$ and $n^{*}$, where $l^{*}(M)=l\left(M^{-1}\right)$ and $n^{*}(M)=n\left(M^{-1}\right)$ for divisors $M$ of $\widetilde{A}$.

The theorem then reads

$$
\begin{equation*}
l^{*}(M)=l^{*}\left(\mathfrak{D}_{M} M^{-1}\right)+n^{*}(M)-r^{2} \xi\left(S_{0}\right) . \tag{4.5}
\end{equation*}
$$

Of course, we could have simply reasserted Theorem 3.4 for $V=A$, but the advantages of restricting our divisors are important: The invariant $\xi(A)=$ $r^{2} \xi\left(S_{0}\right)$ arises, and the existence of inverse and canonical divisors allows a nice expression for the annihilator.

Before comparing this theorem with the Riemann-Roch theorem of Witt [16], we shall need to discuss the different, which will lead to some other points of interest. Accordingly we shall defer the comparison to the next section, where we shall also state a generalization of the theorem of Weil [15].

The different. Define the subspace $M$ of $\bar{S}$ as

$$
\begin{equation*}
M=\{a ; a \epsilon \bar{S}, T(a I) \subset \mathfrak{o}\}, \tag{4.6}
\end{equation*}
$$

where $T$ is the reduced trace from $\bar{S}$ to $R$ obtained by extending that from $r \times r S_{0}$ to $K . \quad M$ is obviously an open two-sided $I$-module and is contained in $b^{-1} I^{\prime}$ for a regular $b \in R$ satisfying $b_{\mathcal{D}} \subset$ kernel $\chi$; therefore $M$ is linearly compact. By the proof of Lemma 4.1, $M$ has the form $M=D^{-1} I$ for some regular element $D$ of $\bar{S}$ of the form (4.3); by Lemma 1.2, each $P$-component of $D$ in the form (4.3) has $a_{P} \in \mathcal{O}_{P}$. The divisor $D I$ is called the different of $I / K$. If $C^{-1} I C$ is any maximal open linearly compact subring of $\bar{S}$, then its different (over $K$ ) is defined analogously and turns out to be $C^{-1} D I C$. The different $D_{x}$ of $K / F(x)$ is defined ${ }^{10}$ similarly by the relation

$$
\left\{a ; a \in R, T_{K / F(x)}(a \mathfrak{o}) \subset \mathfrak{o}_{0}\right\}=D_{x}^{-1} \mathfrak{D}
$$

where $\mathrm{o}_{0}$ is the ring of integers of $R_{0}$. The different $D^{*}$ of $I / F(x)$ is defined similarly, and the result that $D^{*} I=D D_{x} I$ follows immediately from the factorability of the reduced trace.

Now let $\chi$ be any admissible character of $R$, and let $d o$ and $\mathfrak{D} I$ be the corresponding canonical divisors of $R$ and $I$, respectively. We shall prove

Lemma 4.4. $\quad \mathfrak{D} I=D d I$.
Proof. Taking $M$ as in (4.6), we have $(d D)^{-1} I=d^{-1} M$ and $\chi T\left(d^{-1} M I\right)=$ $\chi\left(d^{-1} T(M I)\right)=\chi\left(d^{-1} \mathfrak{v}\right)=0$. Therefore $\mathfrak{D} I \subset d D I$. Since $T$ is open, $T\left(D^{-1} I\right)$ is a divisor of $R$, which must then be $\mathfrak{o}$. Therefore $T\left(d^{-1} D^{-1} I\right)=\mathfrak{o}^{\prime}$. Now $\chi T\left(\mathfrak{b}^{-1} I\right)=0$ implies $T\left(\mathfrak{D}^{-1} I\right) \subset \mathfrak{o}^{\prime}$, or $T\left(d \mathfrak{b}^{-1} I\right) \subset \mathfrak{o}$. Thus $d \mathfrak{D}^{-1} I \subset$ $D^{-1} I$, or $\lesssim I \supset d D I$, Q.E.D.

We can apply this result to the separable extension $K / F(x)$, achieving the result that $d \mathfrak{0}=D_{x} d_{0} \mathfrak{o}$, where $d_{0} \mathfrak{0}_{0}$ is the canonical divisor of $F(x)$ corresponding to the admissible character $\chi_{0}$ of $R_{0}$, and $d \mathrm{o}$ is the canonical divisor of $R$ arising from $\chi=\chi_{0} T_{0}, T_{0}$ being the reduced trace from $K$ to $F(x)$. Furthermore, we can easily construct a $\chi_{0}$ such that $d_{0}=u^{-2}$, where $u$ is the denominator of $x$. (To carry out this construction, one needs the facts that $R_{0}=$

[^5]$F(x)+u \mathrm{D}_{0}$, a topological direct sum, and that $\mathrm{D}_{0}^{\prime} \cap F(x)$ has $F$-dimension equal to the genus of $F(x)$, which is $0,[9,10]$.) Thus we have proved rather simply the well-known result [6, p. 374] that $d=D_{x} / u^{2}$, i.e., that $\left(D_{x} / u^{2}\right)_{\mathfrak{0}}$ is a canonical divisor of $R$. From this result it follows that $\left(D^{*} / u^{2}\right) I$ is a canonical divisor for $I$.

Thus there is an admissible character of $R$ such that if $M=C^{-1} I C_{0}$ is a divisor of $\bar{S}$, then the annihilator of $M$ is the space $C_{0}^{-1}\left(D^{*} / u^{2}\right)^{-1} I C=M^{\prime}$. We have proved

Lemma 4.5. There exists an admissible character $\chi_{1}$ of $R$ such that if $M$ is a divisor of $\widetilde{A}$ with $M_{l}=H$, then the annihilator $M^{\prime}$ corresponding to $\chi_{1}$ is that divisor satisfying

$$
\begin{equation*}
\left(M M^{\prime}\right)^{-1}=u^{-2} \delta(H) \tag{4.7}
\end{equation*}
$$

where $\delta(H)$ is the different of $H / F(x)$.
Since the various differents of $\widetilde{A} / K$ are determined up to conjugation, and multiplication by units, and since $D_{P} \in S_{P}$ for each $P$, the integer $\delta_{P}=v_{P}\left(D_{P}\right)$ is an invariant of $\widetilde{A}_{P}$; in fact, $\delta_{P}$ depends only on $S_{0}$, so that we may as well take $A=S_{0}$ when investigating this differential exponent $\delta_{P}$.

We do so now; our first result will be a relation between the reduced traces $T_{P}$ and $T_{1 P}$ of $S_{P} / K_{P}$ and $\Im_{P} / \Omega_{P}$, respectively. Let $\left[S_{P}: K_{P}\right]=m_{P}^{\prime 2}=$ $e_{P} f_{P}$; thus if $[S: K]=m^{2}$, then $m^{2}=m_{P}^{\prime 2} \mu_{P}^{2}$. For each $a \in \mathcal{O}_{P}$ let $a^{*}$ denote the residue class of $a$ modulo $\mathfrak{P}_{P}$.

Lemma 4.6. For each $P \in \mathfrak{M}$, the reduced traces satisfy

$$
T_{P}(a)^{*}=\alpha_{P} T_{1 P}\left(a^{*}\right), \quad a \in \mathcal{O}_{P}
$$

where $\alpha_{P}=e_{P} m_{P} / m_{P}^{\prime}$ is a positive rational integer.
Proof. The matrix of the regular representation of a generic element of $S_{P}$ (with respect to the integral basis of Lemma 1.2), on reduction modulo $\mathfrak{B}_{P}$ breaks naturally into submatrices of size $f_{P} \times f_{P}$; below the diagonal block these matrices are zero; the diagonal block is an $e_{P}$-fold repetition of the matrix of the regular representation of a generic element of $\Im_{P} / \Omega_{P}$. Therefore the characteristic polynomial $c(X)$ of $S_{P} / K_{P}$ is related to that, $c_{1}\left(X^{*}\right)$, of $\mathfrak{S}_{P} / \Omega_{P}$ by the formula

$$
\begin{equation*}
c(X)^{*}=c_{1}\left(X^{*}\right)^{e_{P}} \tag{4.8}
\end{equation*}
$$

Let $f(X)$ be the minimal polynomial of $S_{P} / K_{P}$ and $f_{1}\left(X^{*}\right)$ that of $\Im_{P} / \Omega_{P}$. We know that in general [1, p. 17] the characteristic polynomial of a skewfield is a power of the minimal polynomial, and that, when the center is separable over the base field, this power is the index of the skew-field [1, p. 123]. Therefore (4.8) becomes

$$
\begin{equation*}
f(X)^{*^{m_{P}^{\prime}}}=f_{1}\left(X^{*}\right)^{m_{P e_{P}}} \tag{4.9}
\end{equation*}
$$

where $m_{P} \geqq 1$ is the index of $\mathfrak{S}_{P}$ whenever the center of $\Im_{P}$ is separable over $\Omega_{P}$. But $f_{1}\left(X^{*}\right)$ is irreducible over $\Omega_{P}\left(\xi_{i}\right)$ (the $\xi_{i}$ being indeterminates used to make the generic element); therefore $f^{*}$ is a power of $f_{1}$ : There is a rational integer $\alpha_{P}$ such that $f(X)^{*}=f_{1}\left(X^{*}\right)^{\alpha_{P}}$, which gives our desired formula. It follows from (4.9) that $\alpha_{P}=e_{P} m_{P} / m_{P}^{\prime}$.

Using Lemma 4.6, the formula of which is analogous to the one for fields, in which $e_{P}$ replaces our $\alpha_{P}$, one proves easily, exactly as in [17, p. 70],

Proposition 4.2. The differential exponent $\delta_{P}$ is at least $e_{P}-1$ for all $P \in \mathfrak{M}$, and $\delta_{P}=e_{P}-1$ if and only if both (i) and (ii) hold:
(i) The center of $\mathfrak{S}_{P}$ is separable over $\Omega_{P}$.
(ii) The characteristic of $F$ does not divide $\alpha_{P}$.

As we shall see in a moment, the condition (i) implies condition (ii) in the noncommutative case, whereas (ii) is necessary for field extensions of $K$. But the relative simplicity of the proof here and the interesting corollaries which obtain at this point may justify the redundancy.

Corollary 4.3. $\quad e_{P}=1$ for almost all $P \in \mathfrak{M}$.
Proof. $\quad \delta_{P}=0$ for almost all $P$, and $\delta_{P} \geqq e_{P}-1 \geqq 0$.
Corollary 4.4 If the constant field $F$ has characteristic zero or is finite, then $\delta_{P}=e_{P}-1$ for all $P \in \mathfrak{M}$. When $F$ is finite, then $e_{P}=f_{P}=m_{P}^{\prime}$, and $\alpha_{P}=1$, both for all $P \in \mathfrak{M}$, so that $K_{P}$ splits $S_{0}$ for almost all $P$.

Proof. The result in characteristic zero is immediate. When $F$ is finite, then $S_{P}$ is a finite skew-field and is therefore a field, by a famous theorem of Wedderburn. Then $m_{P}=1$ for all $P$, so that $\alpha_{P}=e_{P} / m_{P}^{\prime}$. But $e_{P} \mid m_{P}^{\prime}$ in general, by Lemma 1.3; therefore $e_{P}=m_{P}^{\prime}$, or $\alpha_{P}=1$ for all $P$. From $e_{P} f_{P}=m_{P}^{\prime 2}$ follows now $e_{P}=f_{P}=m_{P}^{\prime}$.

We now prove that the assumption on $\alpha_{P}$ in the proposition is unnecessary. The proof is taken essentially from [12, p. 148].

Theorem 4.2. The center of the residue-class skew-field $\mathfrak{S}_{P}$ is separable over $\Omega_{P}$ if and only if the differential exponent $\delta_{P}$ is $e_{P}-1$.

Proof. We need only prove that separability implies $\delta_{P}=e_{P}-1$. It suffices to prove the existence of some $b \in \mathcal{\vartheta}_{P}$ such that $T_{P}(b)$ is a unit in $K_{P}$. To this end let $K_{1}$ be a separable, unramified, maximal subfield of $S_{P} / K_{P}$ with separable residue-class field $\Omega_{1}$ over $\Omega_{P}$ [18, p. 12], [12, p. 148]. Then our formula of Lemma 4.6 for the unramified field extension $K_{1} / K_{P}$ becomes $T_{2}(a)^{*}=T_{21}\left(a^{*}\right), a \in \mathcal{O}_{P} \cap K_{1}$, where $T_{2}$ is the trace of $K_{1} / K_{P}$ and $T_{21}$ that of $\Omega_{1} / \Omega_{P}$. The separability of the last named extension implies that for some $b \in \mathcal{O}_{P} \cap K_{1}, T_{2}(b)$ is a unit of $K_{P}$. But $T_{P}$, when restricted to the maximal subfield $K_{1}$, equals $T_{2}$, Q.E.D.

Note. The existence of $K_{1}$ can be proved directly as follows: Let $L_{1}$ be a separable maximal subfield of $\mathfrak{S}_{P}$ over the center $C_{P}$ of $\mathfrak{S}_{P}[1, \mathrm{p} .57]$.

Then $L_{1}=\Omega_{P}\left(a^{*}\right)$ for some $a \in \mathcal{O}_{P}$, by the separability of $C_{P} / \Omega_{P}$. Consider now $K_{P}(a)$, which has $L_{1}$ as residue-class field. Let

$$
f(X)=X^{q}+a_{1} X^{q-1}+\cdots+a_{q}
$$

be a polynomial over $\mathfrak{o}_{P}$ such that $f(X)^{*}$ is the (irreducible) minimal polynomial of $a^{*}$ over $\Omega_{P}$. Then $f(X)$ is irreducible over $K_{P}$. By Hensel's lemma, $f(X)$ has a zero $b$ in $K_{P}(a) \cap \mathcal{O}_{P}$ such that $b^{*}=a^{*}$. Therefore $K_{P}(b)$ has $\Omega_{P}\left(a^{*}\right)$ as residue-class field, which implies that $K_{P}(b) / K_{P}$ is unramified, being of degree equal to that of $\Omega_{P}\left(a^{*}\right) / \Omega_{P}$. Now we prove that $K_{P}(b)$ is a maximal subfield of $S_{P}$ and therefore that $K_{P}(a)=K_{P}(b)$. If $K_{P}(b)$ were not maximal, then its commutator algebra $K^{\prime}$ in $S_{P}$ would not be a field [1, p. 53]. Then we would have $f\left(K^{\prime} / K_{P}(b)\right)>1$, by Lemma 1.3. This would imply the existence of a proper field extension of $L_{1}$, in contradiction to the maximality of $L_{1}$.

The Riemann-Hurwitz formula. We shall give the analogue of the RiemannHurwitz (? or Zeuthen-Halphen) formula relating, in our case, $\xi\left(S_{0}\right)$ to the genus $g$ of $K$. If $\mathfrak{D O}$ is a canonical divisor of $\bar{S}_{0}$, we know that $n(N(\mathfrak{D}))=$ $2 \xi\left(S_{0}\right)=n(N(D))+n(N(d))$, the latter from Lemma 4.4. Since $d \in R$, $N(d)=d^{m^{2}}$, where $\left[S_{0}: K\right]=m^{2}$. If $D_{P}=\left(\delta_{i j} a_{P}\right), a_{P} \in S_{P}$, the matrix being of size $\mu_{P} \times \mu_{P}$, then, at $P, N(D)=N_{s_{P} / K_{P}}\left(a_{P}\right)^{\mu_{P}^{2}}$. Thus $\nu_{P} N(D)=\mu_{P}^{2} f_{P} \delta_{P}$, where $\delta_{P}$ is the differential exponent. Our above equation now becomes

$$
2 \xi\left(S_{0}\right)=m^{2}(2 g-2)+\sum_{P \in \Re} \mu_{P}^{2} f_{P} n_{P} \delta_{P}
$$

Since $m^{2}=\mu_{P}^{2} e_{P} f_{P}$, we can put this as
Theorem 4.3. Let $S_{0}$ be a skew-field with finite rank $m^{2}$ over the center $K$. Then the invariants $\xi\left(S_{0}\right)$ and $\xi(K)=g-1$ are related by the formula

$$
2 \xi\left(S_{0}\right)=m^{2}\left(2 g-2+\sum_{P \epsilon \mathfrak{M}} n_{P} \delta_{P} / e_{P}\right)
$$

where $\delta_{P}$ is the differential exponent at $P$, discussed in the preceding section.
The invariant $\xi\left(S_{0}\right)$. An interesting question concerning $\xi\left(S_{0}\right)$ is whether the two terms defining it are themselves invariants. That is, does $l\left(H_{1}\right)$ equal $l\left(H_{2}\right)$ for all $H_{1}, H_{2} \in \mathscr{L}\left(\widetilde{S}_{0}\right)$ ? The equivalent question, of course, is whether the same holds with $H_{1}^{\prime}, H_{2}^{\prime}$ in place of $H_{1}, H_{2}$. Some partial results in the affirmative are contained in the following two lemmas. They are phrased in the matrix terminology.

Lemma 4.7. Let $S_{0}$ be a skew-field extension of $K$ with no "constant part"; that is, assume that $F$ is (relatively) algebraically closed in $S_{0}$. Then $l\left(b^{-1} \mathcal{O}\right)=1$ for all regular $b \in \bar{S}_{0}$.

Proof. Let $a \in S_{0}$ and $a \in b^{-1} \mathcal{O b}$. We wish to prove $a \in F$. The characteristic polynomial $c(X ; u)$ of $S_{0} / K$ is the same as that of $K_{P} \otimes_{K} S_{0} / K_{P}$. When specialized to the element $a$, this polynomial has coefficients in $K$. As an element of $b_{P}^{-1}\left(\mu_{P} \times \mu_{P} \mathcal{O}_{P}\right) b_{P}, a$ has the same characteristic polynomial
as $b_{P} a b_{P}^{-1} \in \mu_{P} \times \mu_{P} \mathcal{O}_{P}$ has; this one has coefficients in $\mathfrak{o}_{P}$, however, since, with respect to the usual integral basis $\left\{w_{i} t^{j} \cdot\right.$ (matrix units) $\}$, $b_{P} a b_{P}^{-1}$ has coordinates in $\mathfrak{o}_{P}$. Therefore, the coefficients of $c(X ; a)$ are in $K \cap \mathfrak{D}=F$. By assumption, $a \in F$, Q.E.D.

Corollary (of the proof). If $u_{1}, \cdots, u_{n}$ is a basis of $S_{0} / K$ and $u=x_{1} u_{1}+\cdots+x_{n} u_{n}$ a generic element of $S_{0}$, then the characteristic polynomial of $S_{0} / K$ has coefficients in $F\left[x_{1}, \cdots, x_{n}\right]$.

Our next result shows that in some skew-field extensions obtained entirely by extensions of the constant field, the quantities in question are also invariants.

Lemma 4.8. Let $K$ be a function field of genus 0 , and let $S_{0}$ be a normal skewfield of finite rank $m^{2}$ over $K$. If the different $D \mathcal{O}$ of $\bar{S}_{0}$ satisfies $n(N(D))<2 m^{2}$, then $l\left(a^{-1} \mathcal{O}^{\prime} a\right)=0$ for all regular $a \in \bar{S}_{0}$.

Proof. The relation $y \in S_{0} \cap a^{-1} \mathcal{O}^{\prime} a$ implies $N(y) \in N\left(\delta^{-1}\right) \mathfrak{0}$, where $\mathcal{O}^{\prime}=\mathfrak{D}^{-1} \mathcal{O}$. Using Lemma 4.4, we find that $N\left(\mathfrak{b}^{-1}\right)=d^{-m^{2}} N\left(D^{-1}\right)$. Now $n\left(N\left(b^{-1}\right)\right)=2 m^{2}-n(N(D))>0$ under our assumptions. Therefore $y=0$.

This condition holds, for example, in the case $K=F(x), F$ the field of real numbers, $S_{0}=K(i, j, k)$, where $i^{2}=j^{2}=k^{2}=-1, i j=k, j k=i$, and $k i=j$. Here all $\delta_{P}=0$. It may be of interest to observe that when $P$ has degree $n_{P}=1$, then $f_{P}=4, e_{P}=1$ (hence $\delta_{P}=0$ ), and $S_{P}=K_{P}(i, j, k)$. And when $n_{P}=2$, then $e_{P}=f_{P}=1$ (again $\delta_{P}=0$ ), and $S_{P}=K_{P}$. Thus when $F$ is infinite, it can happen that splitting occurs at infinitely many $P$, and nonsplitting occurs at infinitely many $P$.

## 5. The theorems of Witt and Weil

In comparing Theorem 4.1 with the Riemann-Roch theorem of Witt [16, p. 22] we shall first show that Witt's class of divisors is the same as ours. Witt defines divisors as follows: For a separating element $x$ of $K$, first consider an ideal $M_{0}$ in $A$ with respect to $F[x]$. A finite prime divisor $P$ being one for which $\nu_{P}(x) \geqq 0$, consider the closure $M_{P}$ of $M_{0}$ in $\widetilde{A}_{P}$ for finite $P$. ( $M_{P}$ is shown to be the $P$-component of one of our divisors in Theorem 2.4.) At the (finite number of) nonfinite $P \in \mathfrak{M}$, introduce "components" formally in any possible way. These "components" and the $M_{P}$ define a Witt-divisor. Although Witt does not explicitly define these "components", they can only be normal ideals (i.e., those belonging to maximal orders) in $\widetilde{A}_{P}$ with respect to $0_{P}, P$ nonfinite; otherwise his Satz 3, which says that his class of divisors is independent of $x$, would be false. But our Lemmas 2.4 and 4.1 (local form), plus the obvious fact that a local ideal is open and linearly compact, show that such ideals are $P$-components of our divisors. Therefore, the set of $P$-components $M_{P}$, one for each $P \in \mathfrak{M}$, defining a Witt-divisor $M$, is precisely the set of $P$-components of one of our divisors, and conversely.

Witt defines the degree of a divisor $M$ as the degree in $R_{0}$ of $N_{0}(M)$, where
$N_{0}$ is the norm from $\tilde{A}$ to $R_{0}$. Putting aside the question whether at $P$ this norm maps local maximal orders onto $\mathfrak{D}_{Q}$, let us take the Witt-degree of $M$ as the degree in $R_{0}$ of $N_{0}(M) \mathfrak{D}_{0}$. If $a$ is a regular element of $R$, then the degree $n(a)$ of $a$ in $R$ is the degree $n_{0}\left(N_{*}(a)\right)$ of the norm (from $R$ to $R_{0}$ ) of $a$. Therefore, by our definition in $\S 4$ and our formula (4.2), the Witt-degree of $M$, which we shall denote as $n^{*}(M)$, is the negative of our degree $n(M)$.

Witt defines the quantity $\{M\}$, for a divisor $M$, as $l\left(M^{-1}\right)$, or $l^{*}(M)$, in the notation of our Remark 4.1.

The "complementary" divisor $M^{*}$ to the divisor $M^{-1}$ is defined by Witt ${ }^{11}$ as that satisfying $M^{*} M^{-1}=u^{-2} \delta\left(M_{l}\right)$, in the notation of (4.7). By Lemma $4.5, M^{*-1}$ is what we call $M^{\prime}$.

Witt defines a genus $G$ of $A$ by the formula $2 G-2=n^{*}\left(u^{-2} \delta\left(M_{l}\right)\right) . \quad$ By Corollary 4.1, $n^{*}\left(u^{-2} \delta\left(M_{l}\right)\right)=2 r^{2} \xi\left(S_{0}\right)$.

Finally, Witt states his Riemann-Roch theorem as

$$
l^{*}(M)=l^{*}\left(M^{*}\right)+n^{*}(M)-G+1
$$

which agrees with our (4.5) when we determine the $\delta_{M}$ there by means of the admissible character $\chi_{1}$ of Lemma 4.5.

We shall now sketch the proof of a theorem which includes that of Weil [15, Ch. I, 3] (with trivial signature). Let $V=r \times r^{\prime} S_{0}$ be the space of all matrices of $r$ rows and $r^{\prime}$ columns over $S_{0}$, a skew-field of finite rank over the center $K$. As $A$ and $B$ we take $r \times r S_{0}$ and $r^{\prime} \times r^{\prime} S_{0}$, the actions being the usual matrix multiplication. Let $T$ denote the reduced trace from $S_{0}$ to $K$ and $\operatorname{Tr}$ the ordinary matrix trace. Then for $V^{\prime}$ we take $r^{\prime} \times r S_{0}$, and we set $\left\langle v, v^{\prime}\right\rangle_{0}=T \operatorname{Tr}\left(v v^{\prime}\right), v \in V, v^{\prime} \in V^{\prime}$. Our dual pairing of $\widetilde{V}$ and $\tilde{V}^{\prime}$ to $F$ is then $\left[\tilde{v}, \tilde{v}^{\prime}\right]=\chi T \operatorname{Tr}\left(\tilde{v} \tilde{v}^{\prime}\right)$, for $\tilde{v} \in \tilde{V}, \tilde{v}^{\prime} \in \tilde{V}^{\prime}$.

The numbers $\rho_{1}$ and $\rho_{2}$ of Lemma 3.3 are $r^{\prime} / r$ and $r / r^{\prime}$, respectively.
Letting $H$ denote a member of $\mathscr{L}\left(S_{0}\right)$, we take for $V$-divisors subspaces of $\tilde{V}$ of the form $M=a^{-1}\left(r \times r^{\prime} H\right) b, a$ regular in $\widetilde{A}, b$ regular in $\widetilde{B}$ ( $V^{\prime}$-divisors have the form $\left.b^{-1}\left(r^{\prime} \times r H\right) a\right)$. Our annihilator $M^{\prime}$ is $b^{-1}\left(r^{\prime} \times r b^{-1} H\right) a$. The degree (3.5) of a divisor $M$ is invariant with respect to all divisors of the form $c^{-1}\left(r \times r^{\prime} H\right) c, c$ regular in $\widetilde{S}_{0}$. We define the degree of $M$ as

$$
\nu\left(M, r \times r^{\prime} H\right)=n\left(N(a)^{r^{\prime} / r} N^{\prime}(b)^{r / r^{\prime}}\right)
$$

where $N$ is the norm from $\widetilde{A}$ to $R$ and $N^{\prime}$ that from $\widetilde{B}$ to $R$.
The quantity $l\left(c^{-1}\left(r \times r^{\prime} H^{\prime}\right) c\right)-l\left(c^{-1}\left(r \times r^{\prime} H\right) c\right)$ is the same for all regular $c \in \widetilde{S}_{0}$; it equals $r r^{\prime} \xi\left(S_{0}\right)$.

Our generalization of the Riemann-Roch theorem to $V$ is
Theorem 5.1. If $M$ is a $V$-divisor, then

$$
l(M)=l\left(M^{\prime}\right)+n(M)-r r^{\prime} \xi\left(S_{0}\right) .
$$

[^6]This theorem agrees with that of Weil when his signature is identically 1 if we take $S_{0}$ to be $K$ and $F$ the field of complex numbers. To verify this agreement, one uses the one-one correspondence between canonical divisors and differentials of $K$ (see [9]).

We could derive Theorem 4.1 from Theorem 5.1 if we used the isomorphism of Theorem 2.3.

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[^1]:    ${ }^{2}$ By the term "function field," we mean a field $K$, containing a subfield $F$, (relatively) algebraically closed in $K$, and containing an element $x$ not in $F$, such that $K / F(x)$ has finite degree.
    ${ }^{3}$ A vector space over $F$ is linearly topologized if its additive group is a topological group for which some collection of linear subspaces serves as a fundamental system of neighborhoods of 0 . For the theory of such spaces, including the duality theorem, see Lefschetz [11, pp. 72-83]. The analogue of the "second isomorphism theorem" for these spaces is proved in [9].
    ${ }^{4}$ For the structure theory of algebras assumed in this paper, see, for example, [2].

[^2]:    ${ }^{5}$ See, for example, the proof of Theorem 5, p. 61, of [17].
    ${ }^{6}$ For the definition and properties of the reduced trace, see [1, pp. 122-125].
    ${ }^{7} \mathrm{An}$ account of orders and their ideals will be found in [4, Ch. VI].

[^3]:    ${ }^{8}$ That (2.1) holds even when $x$ is not a separating element of $K$ was proved by Iwasawa in [10, §3].

[^4]:    ${ }^{9}$ For proofs of these properties of the $\nu$-function, see $[10, \S 1]$.

[^5]:    ${ }^{10}$ Here we need to assume, with Witt, that $K / F(x)$ is separable.

[^6]:    ${ }^{11}$ The different appearing in Witt's definition is misprinted as that of $K / F(x)$ instead of $M_{l} / F(x)$.

