Continuum Mechanics By Donald E. Carlson 1991

These notes are intended to be a textbook for an introductory but careful study of modern continuum mechanics at the advanced undergraduate on beginning graduate level. They are directed at American engineering students who have been exposed to the mathematics courses (calculus of one and several val variables, matrix algebra, ordinary differential equations) and the engineering mechanics courses (statics and dynamics of particles and rigid bodies, elementary mechanics of fluids and deformable solids) of the traditional curriculum. However, with respect to mechanics, the notes are completely self-contained; and consequently, they should be accessible to students of mathematics and the sciences.

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Part A

Mathematical Preliminaries

This long first part provides the additional mathematical tackground record for our study. Some of the points of view adopted
here are more common in mathematics than in engineering, but
they are particularly useful in continuum mechanics. Keeping the
mathematical preliminaries separate from the devlopment of the
nechanics has the advantage of permitting a cleaner treatment,
and it also serves to emphasize that these mathematical
notions are very general and are applicable to areas other
than continuum mechanics. However, so as not to unotally
belay reaching continuum mechanics proper, the study of some
of the sections could be postponed until just before they are
needed. E.g., Curilinear condinate systems are not employed
until the development of particular solutions in Part C.
Such matters will be pointed out at the appropriate places
in the text.

Chapter A1

Some Basic Concepts

A1.1. Sets and Lists

A set is a collection of objects viewed as a single entity. A sushel of apples, a head of cows, the positive integers, the real numbers, the exercises of Calculus III are all examples of sets, The objects in the collection are called the elements or members of the set, and the notation

$a \in A$

is used to indicate that a is an element of set A, E.g., $\frac{3}{2}z$ is an element of the set of real numbers R, and we write

3/2 ∈R,

It is important to realize that the above remarks to not constitute an adequate definition of set unless we know the meaning of collection and object. Here, we follow the tradition of taking the notion of set to be a primitive concept; i.e., an undefined but intuitively natural concept. Collection and object are just everyday nontechnical words that help us elucidate the primitive concept of set. Once primitive notions have fell laid down, then mathematics customarily proceeds

more formally through precise definitions and therems in terms of the primitive concepts. This will be our approach throughout these rotes,

Definition A1.1.1. Two sets A and B are said to be equal, and we write A = B, if they contain the same elements,

Here, and subsequently, definitions are understood as "if and only if" statements. Thus, the above definition says that if the sets A and B contain the same elements, then A=B, and conversely if A=B, then A and B contain the same elements,

Definition A1.1.2, \underline{A} set \underline{A} is said to be a subset of a set \underline{B} , and we write $\underline{A} \subset \underline{B}$ on $\underline{B} \supset A$, if every element of \underline{A} is also an element of \underline{B} ; i.e., if

 $a \in A \Rightarrow a \in B$,

If A < B but A ≠ B, then A is a proper subset of B.

As usual, the symbol \Rightarrow stands for "implies", and the slash / is used for regation; thus, $A \neq B$ means that the sets A and B are not equal.

The following theorem provides the standard tool for proving that two sets are equal.

Theorem A1.1.1, Let A and B be sets, Then A=Biff
both A < B and B < A.

Here, if is an abbreviation for "if and only if", Thus, the theorem really consists in the two statements:

(i) A=B if both A < B and B < A;

(ii) A=B only if both ACB and BCA.

The first statement is straightforward, although, at least for me, it reads more smoothly in the form

(i) If both A < Band B < A, then A = B.

The seemed statement is more subtle. It means that the touth of A = B requires that both A < B and B < A. Thus, in phrasing parallel to that used in (i), statement (ii) can be cast as

(ii) If A=B, then both ACB and BCA.

In the terminology of elementary logic, statement (ii) is the converse of statement (i), and vice versa,

Common alternative ways of stating Theorem A1.1.1 would

A=B is equivalent to both A = B and B < A;

Necessary and sufficient for A = B is that $A \subset B$ and $B \subset A$, $A = B \iff B$ oth $A \subset B$ and $B \subset A$.

If course, the symbol \Leftrightarrow is a generalization of \Rightarrow , it is read as "implies and is implied by".

Since Theorem A1.1.1 provides a condition which is equivolent to A=B, it is often taken as the definition of A=B. In such an approach, our Definition A1.1.1 would then be a theorem. We shall always attempt to adopt the definitions that are the most natural rather than those which lead to the most efficient development. Of course, considerable subjectivity is involved in such Secisions.

Before going on we must not forget to tuen to the

Proof of Theorem A1.1.1. Suppose that both ACB and BCA. Since ACB, it follows from Definition A1.1.2 that every element of A is also in B, but since BCA, every element of B is also in A. Thus, there are no elements of B that are not also in A. In other words, A and B contain exactly the same elements. Therefore, A = B by Definition A1.1.1.

Conversely, suppose that A = B. Consider a typical element

 $a \in A$. Since A = B, A and B contain the same elements. Therefore, we also have $a \in B$. Thus, $a \in A \Rightarrow a \in B$; i.e., $A \subset B$ by Definition A1.1.2, In the same manner, we see that $B \subset A$. Hence, $A = B \Rightarrow b$ oth $A \subset B$ and $B \subset A$. \Box

The symbol is used here to indicate that the end of a proof has been reached.

Sets which contain just a few clements are usually indicated by listing all of their clements between between between between

{a,b,c},

No significance is attached to the order in which the elements are listed, a set which contains a single element, say $\{a\}$, is called a singleton.

In more complicated cases, particular sets are specified by their elements possessing some refining property. The notation

 $\{x:P(x)\}$

stands for the set of all elements & which have the property

^{&#}x27;gfen, {x | P(x)}.

Par. E.g.,

{x: x e R, 1x1<1}

is the set of all real numbers with absolute value loss than 1. Sometimes it is advantageous to delimit the possibilities for candidacy in a set at the outset. The notation

{ oc ∈ A: P(x)}

stands for the set which consists in all of those elements a in A which passess the property P(x). Thus, our previous example could have been expressed as

{x∈R: |x|<1}.

Next we tuen to the common operations with sets.

Definition A1.1.3. The union of two sets A and B is the set $AUB = \{ x : x \in A \text{ or } x \in B \}$.

We will always use the conjunction or in its "weak" sense. Thus, in the above definition, we could have both $x \in A$ and $x \in B$,

Definition Al. 1.4. The intersection of two sets A and B is the set

ANB = {x: x ∈ A and x ∈ B}.

Of course, it could happen that the sets A and B would have no elements in common, then their intersection would be a set with no elements. This passibility leads us to the next two definitions.

empty set! It is denoted by p.

Definition A1.1.b. Let A and B be set, If $A \cap B = \emptyset,$

then A and B are said to be disjoint or nonintersecting.

The next theorem gives the most important properties of the set operations of union and intersection.

Theorem A1. 1.2. Let A, B, and C be sets. Then we have the idempotent properties

(i) AUA = A, ANA = A;

the commutative properties

Often, the void set on the vacuous set.

(ii) AUB = BUA, ANB = BNA;

the associative properties

(iii) AU(BUC) = (AUB)UC, (ANB)NC = AN(BNC);

and the distributive properties

(iv) AU(BNC) = (AUB) N(AUC), AN(BUC)=(ANB) U(ANC).

Proof. Properties (i) and (ii) are so obvious that we do not supply found people of them here.

To establish (iii), , we first suppose that $x \in (AUB)UC$. Then by repeated use Definition A1.1.3, we have $x \in (AUB)$ on $x \in C$, so that $x \in A$ or $x \in B$ or $x \in C$; consequently, $x \in A$ or $x \in (BUC)$, which means that $x \in AU(BUC)$. Therefore, by Definition A1.1.2, $(AUB)UC \subset AU(BUC)$. In the same way, we see that $AU(BUC) \subset (AUB)UC$. Hence, by Theorem A1.1.1, AU(BUC) = (AUB)UC.

The people of (iii) is left as Exercise A1.1.1.

The proofs of the Listributive peoperties we a little more

More precisely, (iv), says that the operation of union is distributive w.r.t. the operation of intersection. Similarly, for (iv)2. Here, w.r.t. stands for "with respect to".

difficult. Let us consider (iv), . Suppose that $x \in AU(BNC)$. Then by Definition A1.1.3, $x \in A$ or $x \in BNC$, so that by Definition A1.1.4, $x \in A$ or $(x \in B)$ and $x \in C)$. Note that the use of the parenthesis or some such verice is essential here if we are to say what we mean, Next we examine the two possibilities at which we have arrived.

If $x \in A$, then it is correct to say that $x \in AUB$. This is not a natural step, because there is a loss of information; monetheless, it is correct. Similarly, $x \in A \Rightarrow x \in AUC$. Putting these implications together, we have that $x \in A \Rightarrow x \in (AUB) \cap (AUC)$.

If $x \in B$ and $x \in C$, then $x \in AUB$ and $x \in AUC$; and these last two results $\Rightarrow x \in (AUB) \cap (AUC)$.

Thus, both of the possibilites under consideration lead to the same conclusion, and we have that

x∈AU(BNC) ⇒ x∈ (AUB) N (AUC),

or

AU(BNC) < (AUB) (AUC)

by Difinition A1.1.2.

Next suppose that $x \in (AUB) \cap (AUC)$. In accordance with Definitions A1.1.3 and 4, this means that $x \in (AUB)$ and $x \in (AUC)$, which is to say that $(x \in A \text{ or } x \in B)$ and

(x ∈ A on x ∈ C). Let us write out all of the possibilities:

 $x \in A \text{ and } x \in A \Rightarrow x \in AU(BAC)$;

 $x \in A$ and $x \in C \Rightarrow x \in AU(BAC)$;

x∈ B and x∈A ⇒ x∈AU(BAC);

XEB and XEC => XEBAC => XEAU(BAC),

again all of the possibilities lead to the same canclusion, and we have that

 $x \in (AUB) \cap (AUC) \Rightarrow x \in AU(B\cap C)$,

or

(AUB) N (AUC) - AU (BNC)

By Definition A1.1.2. Therefore, by Theorem A1.1.1, AU(BNC) = (AUB) N (AUC),

The proof of liv)2 is left as Exercise A1.1.2. [

In view of the associative properties of union and intersection, we can drop the parentheses and write

AUBUC and ANBNC

without serious ambiguity.

In occasion we shall need to remove the elements of one set from another set. This leads to

Definition A1.1.7. The set-difference of two sets A and B is the set

 $A-B = \{x \in A : x \notin B\},$

Of course, $x \notin B$ means that x is not an element of B,

The operation of set-difference has many interesting peoperties in conjunction with the operations of union and intersection, but we shall need only the definition.

In general, a set which consists of the two elements a and b would be denoted by Ea, b}; and since no meaning is assigned to the order of the listing, {b,a} is exactly the same set. However, on occasion, it is useful to order the elements, and this leads us to

Usually, the set-difference A-B is called the complement of B relative to A and is denoted by A\B. Less often, it is denoted by A\B. Complement is restricted to the case where BCA.

the primitive concept of a "list".

Let n be a stuctly positive integer. Then a list of length n) is an object senoted by

 $(a_n, a_2, ..., a_n)$

which consists of a first element a_1 , a second element a_2 , ..., and an nth element a_n .

Definition A1.1.8. Two lists of order n (9, a2, -, an) and (b, b2, -, bn) are equal, and we write

 $(a_1, a_2, \dots, a_n) = (b_1, b_2, \dots, b_n)_g$ $a_i = b_i \ \forall \ i \in \{1, 2, \dots, n\},$

The symbol & is read as "for all" or "for every". If course, $a_i = b_i$ means that a_i and b_i are the same object.

Precisely because of the ordering built into the notion of a list, a list is not a set. Comparison of the definitions of equality makes this especially

¹ Strictly positive means that O is not included.

² Often, n-tuple.

clear; e.g.,

 $\{\pi, e\} = \{e, \pi\}$ but $(\pi, e) \neq (e, \pi)$.

When the length of a list is small', special terminology is customary.

Definition A1, 1, 9. Lists of length 2, 3, and 4 are called ordered pains, ordered triples, and ordered quadruples, respectively.

Exercise A1.1.3. Is a list of length 1 a singleton?

See the book by Halmos included in the Supplementary Reading at the end of this section for an approach to lists which is entirely in terms of sets,

The following definition will play a major role in our treatment of functions in the next section.

Definition A1, 1, 10, The set-product of two sets A and B

[&]quot;you should resist the temptation to its speak of short lists, unless you know of some short integers.

Almost always, Cartesian product; sometimes, direct product. This is the first instance of our general practice.

is the set of ordered pairs

AxB = {(a,b): a ∈ A, b ∈ B}.

Note that if either A on B is empty, then we cannot four any ordered pairs (a,b). Thus, $\phi \times B = A \times \phi = \phi$.

of giving descriptive rather than proper names to concepts and theorems even when the proper name is used almost universally. In this regard, as well as many others, our thinking has been travily in fluenced by the work of WALTER NOLL.

Supplementary Reading

BARTLE, The Elements of Real Analysis

BISHOP and GOLDBERG, Tensor Analysis on Manifolds

BOWEN and WANG, Introduction to Vectors and Tensors, Vol. 1, Linear and Multilinear Algebra

HALMOS, Naive Set Theory

KOLMOGOROV and FOMIN, Introductory Real Analysis

LOOMIS and STERNBERG, Advanced Calculus

MICHEL and HERGET, Mathematical Foundations in Engineering and Science

NAYLOR and SELL, Linear Operator Theory in Engineering and Science

NOLL, Finite-Dimensimal Spaces

ODEN, Applied Functional Analysis

SIMMONS, Introduction to Topology and Modern Analysis

A1, 2. Functions

Beginning students usually think of "functions" as definite formulas relating real numbers, such as

$$y = f(x) = x^2 - 1$$
 for $x \in \mathbb{R}$.

Of course, to limit oneself to real numbers and definite formulas is clearly two restrictive for the purposes of either mathematics or its applications.

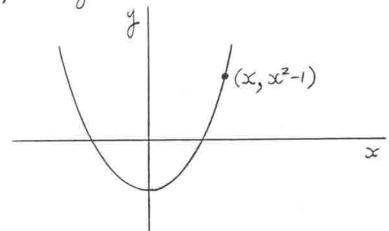
Those who use mathematics typically generalize this along the following lines: "Given two sets A and B, a function + on A to B is a rule of correspondence that assigns to each $x \in A$ a unique element $f(x) \in B$." This definition works quite well, but it has the defect of being based on another primitive notion; namely, "rule of correspondence". Here, we have a dilemma. Obviously, we would be wise to minimize the number of primitive concepts introduced; indeed, it is the current fashim to found all of mathematics on sets, On the other hand, we would like to progress rapidly, perhaps even eventually getting to continuum mechanics, without being build under a tangle of technical Sevelopments. Actually, we were in a similar position at the end of the previous section. There we chose to take the notion of list as a primitive concept, eventhough a set-theoretic definition was available. However, the concept of list is much more transparent than "rule of

comes pondence".

To gain some insight into a set-theoretic approach to "functions", we return to the palestrian example mentioned at the Deginning of the section:

$$y = f(x) = x^2 - 1$$
, $x \in \mathbb{R}$,

The correspondence between x and y can be depicted by the "graph" of x^2-1 .



But the graph is just a special set of ordered pains, and thus we are lead to the following definition. In order to state it concisely, we first introduce the symbols \exists and \ni , which are read as "there exists" and "such that", respectively.

Jan use of > for "such that" Lisallows the use of the symbol ∈ in the backward fashion A>a, which is sometimes read as "the set A contains the element a".

Definition A1, 2.1. Given two sets X and Y, a function I f on X to Y is a subset of XXX with the property that

for each x ∈ X ∃ a unique y ∈ Y ∋ (x,y) ∈ f.

The sets X and Y are called the domain and the codomain of the function f, respectively. For $x \in X$, the unique $y \in Y \ni (x,y) \in f$ is called the value of f at x; it is denoted by

y = f(x).

The uniqueness condition above is conveniently expressed as $(x,y) \in f$ and $(x,y') \in f \Rightarrow y' = y$.

If this requirement is shopped, then we are left with the definition of a relation.

We shall often use the notation $f: X \to Y$

in place of the statement "f is a function on I to I".

Yten, map, mapping, transformation, operator.

Occasionally, the notation $.x \mapsto f(x) = y$

is a useful substitute for

 $(x,y) \in f$ or y = f(x).

Definition A1.2, 2. Let $f: X \to Y$. If $S \subset X$, then the set $f(S) = \{f(x): x \in S\}$

is called the direct image of S under f. The set f(X) is called the range of f. If f(X) = Y, then f is said to be onto Z.

The context will usually prevent the confusion which is possible between function values and direct images.

Let us return to our earlier example of the function $f: \mathbb{R} \to \mathbb{R}$ defined by the rule

 $y = f(x) = x^2 - 1$.

Often, simply image.

²⁰ ften, surjective.

Is this really a function? Here, the domain and coobmain each are R as claimed, indeed, if $x \in \mathbb{R}$, then $y = f(x) = x^2 - 1 \in \mathbb{R}$. To see that the given rule actually defines a function, we must establish the uniqueness condition that

on
$$(x,y) \in f$$
 and $(x,y') \in f \implies y'=y$ on $y = f(x)$ and $y' = f(x) \implies y'=y$ or $y = x^2-1$ and $y'=x^2-1 \implies y'=y$.

Since the meaning of x^2-1 is unambiguous, the last implication obviously is valid; and the rule $f(x) = x^2-1$ refines a function $f: R \rightarrow R$ in accordance with Definition A1.2.1. Of course, the actual function f is the set of ordered pairs $f = \{(x, x^2-1): x \in R\}$.

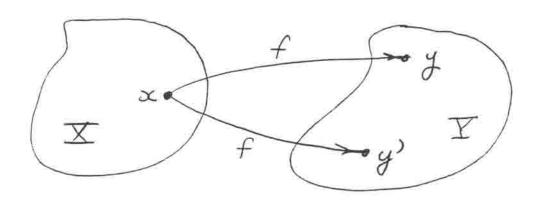
by such rules, we foundize the slove result as

Theorem A1.2.1. An unambiguous evaluation wile, for $x \in X$, $y = f(x) \in Y$,

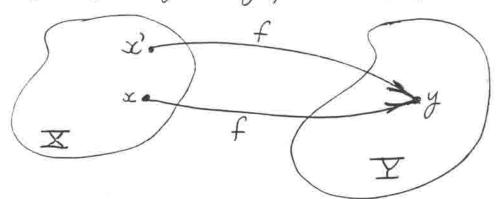
Ufines a function $f: X \rightarrow Y$ through $f = \{(x, f(x)): x \in X\}$

Exercise A1, 2, 1, Is the function $f: \mathbb{R} \to \mathbb{R}$ defined by $f(x) = x^2 - 1$ onto \mathbb{R} ?

By definition, as has been just emphasized, a function $f: X \to Y$ associates to a given element $x \in X$ a unique element $y = f(x) \in Y$, thus, the following picture is not possible.



But this does not rule out the possibility that y is also the volve of f at some other element $x \in X$; i.e., the following picture is possible.



For the important class of "one-to-one" functions, the situation pictured lastly above also is unled out. Formally, we have

Definition A1.2.3. A function $f: X \to Y$ is said to be one-to-one if $\forall x, x \in X$

 $f(x') = f(x) \Longrightarrow x' = x.$

For one-to-one functions, it is possible to " rum the functional correspondence backwards".

Theorem A1.2.2. If $f: X \to Y$ is one-to-one, then the set of ordered pairs

 $f^{-1} := \{ (y,x) \in f(X) \times X : (x,y) \in f \}$

is a function on f(X) to X. Moreover, f - is

¹⁰ sten, one-one, 1-1, injective.

² The symbol := means that the left-hand side is defined by the right-hand side, similarly, for =:.

one-to-one and onto, and

 $f^{-1}(y) = x \iff f(x) = y$.

The function f is said to be the inverse function to the me-to-one function f,

Proof. Refer to Definition A1.2.1. First, we note that by its construction f^{-1} is a subset of $f(X) \times X$, and thus it is a candidate for a function on f(X) to X. We must show that

for each $y \in f(X) \exists a unique x \in X \ni (y,x) \in f^{-1}$.

Accordingly, let $y \in f(I)$. By Definition A1,2,2,

 $f(X) = \{f(x) : x \in X\}$;

and :.) \exists an $x \in X \Rightarrow f(x) = y$. But this equation is just another way of stating that $(x,y) \in f$. Looking back to the definition of f^{-1} , we see that this proves that

for each $y \in f(X) \exists \text{ an } x \in X \ni (y, x) \in f^{-1}$,

Functions that are one-to-one and onto, i.e., injective and surjective, are often said to be bijective.

² As always, the symbol: stands for "therefore".

The issue now is the uniqueness of x, Suppose we have both $(y,x) \in f^{-1}$ and $(y,x') \in f^{-1}$. By the construction of f^{-1} , this means that $(x,y) \in f$ and $(x,y) \in f$ or y = f(x) and y = f(x'). Hence, f(x') = f(x); but f is one-to-one, so by Definition A1.2.3, x' = x. Thus, x is unique, and f^{-1} is a function.

Next we turn to the peoof of $f'(y) = x \iff f(x) = y$.

Now that we know that f' is a function so that the notation f''(y) = x makes sense, the Louble implication above is really tautologous with

 $f^{-1} = \{ (y,x) \in f(X) \times X : (x,y) \in f \}$.

We have $f^{-1}(y) = x \Rightarrow (y,x) \in f^{-1} \Rightarrow (x,y) \in f \Rightarrow y = f(x)$

 $f(x) = y \Rightarrow (x,y) \in f \Rightarrow (y,x) \in f \Rightarrow x = f'(y)$.

To prove that f^{-1} is one-to-one, we must show that $\forall y, y' \in f(X)$

$$f^{-1}(y') = f^{-1}(y) \Rightarrow y' = y$$
.

Accordingly, let y and y' be arbitrary elements of f(X),

and set
$$f^{-1}(y) = x$$
 and $f^{-1}(y') = x'$. Then

$$f^{-1}(y') = f^{-1}(y) \Rightarrow x' = x \Rightarrow f(x') = f(x) \Rightarrow y' = y$$

To prove that f^{-1} : $f(X) \to X$ is onto, we must show in accordance with Definition A1.2.2 that

$$f^{-1}(f(X)) = X,$$

By this same definition

$$f^{-1}(f(X)) = \{f^{-1}(y) : y \in f(X)\}.$$

The problem here is to show that two sets are equal; as usual, we use theorem A1.1.1. Suppose $x \in X$. Set y = f(x). Then $y \in f(X)$ and x = f'(y). Thus,

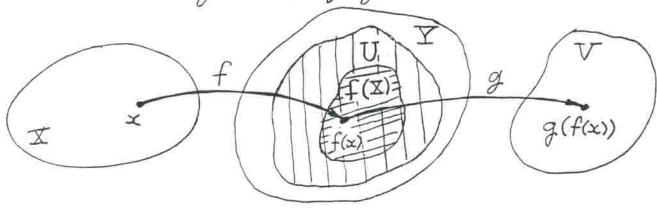
and since x is an arbitrary element of X, $X \subset f^{-1}(f(X))$. Next suppose $x \in f^{-1}(f(X))$. Since $f^{-1}(f(X)) = \{f^{-1}(y) : y \in f(X)\}$, \exists some $y \in f(X) \ni$

 $x = f^{-1}(y) \Rightarrow (y,x) \in f^{-1} = \{(y,x) \in f(X) \times X : (x,y) \in f\} \Rightarrow x \in X,$

so that
$$f^{-1}(f(X)) \subset X$$
. ... $f^{-1}(f(X)) = X$. \Box

Exercise A1.2.2. Show that the function $f: \mathbb{R} \to \mathbb{R}$ defined by $f(x) = x^2 - 1$ is not one-to-one.

Often in continuum mechanics it is necessary to "compose" functions in the following sense, Suppose we have two functions $f: X \to Y$ and $g: U \to V$. If $x \in X$, then f maps x into some element f(x) in the range of f. If f(x) is in the Lomain U of g, then g maps f(x) into some element of V. Thus, we have a scheme for mapping from X to V.



The following theorem makes this procise,

Theorem A1, 2.3. Let $f: X \to Y$ and $g: U \to V$ with $f(X) \subset U$. Then

 $g \circ f := \{(x, v) : x \in X, v = g(f(x))\}$

is a function on X to V, and $(g \circ f)(x) = g(f(x))$.

gof is called the composition of f and g.

Parf. Exercise A1.2.3. [

We shall have need of

Theorem A1.2.4. The composition of functions is associative; i.e., let the domains, codomains, and renges of the functions $f, g, \frac{and}{h}$ be such that the compositions $(h \circ g) \circ f$ and $h \circ (g \circ f)$ make sense, then $(h \circ g) \circ f = h \circ (g \circ f)$.

Beg. Exercise A1.2, 4, 1

Because of Theorem A1, 2, 4, we can write hogof without serious ambiguity.

Exercise A1.2.6. Given a set X, the identity map on X is the function $i_{\underline{x}}: X \to X$ defined by the evaluation rule $i_{\underline{x}}(x) = x$.

For a one-to-one function $f: X \to Y$, show that $f' \circ f = i_X \quad \text{and} \quad f \circ f'' = i_Y.$

Of course, the equality of two functions means that their respective domains and codomains are equal and that the defining subsets of endered pairs (of Definition A1.2.1) are equal. It is left as Exercise A1.2.5 to show that (cont.)

Occasionally, we need to restrict our attention to a subset of the Lomain of a given function, and then the following Lefinition will be aseful.

Definition A1. 2.4, Given a function $f: X \to Y$, let $A \subset X \subset B$. Then the restriction of f to A, denoted by $f|_A$, is the function $f|_A: A \to Y$ defined by the evaluation rule

for $x \in A$, f(x) = f(x)

and a function $g: B \to C$ is an extension of f to B if $g \mid_{X} = f$. Y^{Z}

Same as for & A1.1

^{9.} Theorem A1,2.1,

² This, of course, requires that YCC.

⁽cont. from p. A1.2,13) two functions f and \overline{f} on X to Y are equal iff $\overline{f}(x) = f(x) \ \forall \ x \in X$.

² G. Theorem A1.2,1,

A1.3. Groups

General theories are built up in mathematics by imposing additional structure on sets. The "group structure" is one which occurs repeatedly.

Definition A1.3.1. A group is a set, say I, equipped with a function from Ix I to I, called combination and denoted by

 $(a,b) \longmapsto a*b$,

which satisfies the following axioms:

- (G1) Associativity. $\forall a,b,c \in \mathcal{Y}$ (a*b)*c = a*(b*c);
- (G2) Existence of a neutral element. $\exists n \in \mathcal{S} \Rightarrow$ $n*a = a*n = a \quad \forall a \in \mathcal{S}$,
- (63) Existence of reverse elements. For each $\alpha \in \mathcal{Y} \ni \hat{a} \notin \mathcal{Y} = a \times \hat{a} = n$.

If, in addition, the commutative peoperty

axb=bxa y a, b∈y

is satisfied, then I is said to be a commutative group,

Since the contination function in the definition of a genup may be thought of as taking two elements from I and producing an element also in I, it is often referred to as a closed binary operation. The closure property, i.e., the fact that the function value arb is also in I, is often called the function value arb is also in I, is often called the functional group property.

There are many familiar examples of groups, but they are not always identified as such. The set of real numbers IR is a group w.r.t. acklition. If course, in this cancrete example, we write a+b instead of a*b. I.e., the combination function under consideration is defined on IR *R to R by

 $(a,b) \mapsto a+b$.

Since this rule is unambiguous, the operation of addition in R really is a function by Theorem A1.2.1. What about the three axioms? It is a familiar property of the real numbers that

(a+b)+c = a+(b+c) ∀a,b,c ∈R,

Sometimes, Abelian.

so associativity is satisfied. The number 0 is the neutral element; i.e., $0+\alpha=a+0 \ \forall \ \alpha \in \mathbb{R}$,

The reverse elements are the negatives; i.e., given any $a \in \mathbb{R}$, the number -a has the property

(-a) + a = a + (-a) = 0.

Thus, R is, indeed, a group w.r.t. addition; in fact, since $a+b=b+a \ \forall \ a,b\in R$,

it is a commutative group.

Exercise A1.3.1. Show that $R - \{0\} = : \mathbb{R}^{\times}$ is a commutative group w.r.t. multiplication.

The following result is often useful in manipulations in a group. It is so transparent that usually it is not stated formally. It simply says that if equals are contined with equals the result are equal. We state it here in the language of Definition A1.3.1.

Therem A1.3.1. Let a, b, c, and d be elements of a group,

Then $a = b \text{ and } c = d \implies a * c = b * d.$

Proof. We start with the identity

$$a * c = a * c$$

= $b * c$ (substitution of $a = b$)
= $b * d$ (substitution of $c = d$). \square

Our next theorem asserts that the equation $a \times x = b$ has the unique solution $x = \hat{a} \times b$, with a similar result for $y \neq a = b$.

Theorem A1.3.2. Let \mathcal{Y} be a group. Then for any $a,b\in\mathcal{Y}$, \exists a unique $x\in\mathcal{Y}$ \Rightarrow

 $a \times x = 6$.

In fact,

 $x = \hat{a} * b$,

Similarly, I a unique y ∈ I >

y*a=b,

 $\underline{In\ fact}$, $y = b * \hat{a}$.

Ray. Suppose, tentatively, that \exists an $x \in S \ni a * x = b$. Then combining on the left with \ddot{a} , we have

~ (a*x) = ~ b,

in accordance with Theorem A1.3.1. Now

$$a \times (a \times x) = (a \times a) \times x$$
 (associativity)
= $n \times x$ (reverse)
= $x \times (new + x + x)$,

and i. $x = \delta \times b$. Thus, if x exists, if must be precisely $\delta \times b$. This establishes the uniqueness. To prove existence, we must show that this x has the desired property. First of all, we note that by closure $\delta \times b \in \mathcal{I}$. Next, consider

a*(a*b) = (a*a)*b (associativity)= n*b (reverse) = b (rentral),

Hence, $x = \tilde{a} * b$, does the j'ob.

The "y-problem" can be handled in the same way. Work this out as Exercise A1.3, 7. []

The next five results are all corollaries of Theorem A1.3.2. They also can be viewed as siret consequences of Theorem A1.3.1. We continue to use the language of Definition A1.3.1.

Theorem A1.3.3. (Cancellation Property) Let \mathcal{G} be a group. Then for $a,b,c\in\mathcal{G}$

b*a = c*a ⇒ b=c

and

 $a*b = a*c \implies b = c$,

Proof. Suppose bx a = c * a. We apply the second part of theorem A1.3.2 with b playing the role of the unborning and cx a being the right-hand side b to get

$$b = (c*a)*a$$

$$= c*(a*a) (associativity)$$

$$= c*n (reverse)$$

$$= c (reutral),$$

Otviously, the second assertion can be hardled in the same fastion. But, for Liversity, we leave it for the student as Exercise A1.3.3 to establish if by combining a + b = a + c on the left with a in accordance with Theorem A1.3.1. \square

Theorem A1,3.4. A given group has only one rentral

Proof. Suppose n and n' are both newtral elements for a group \mathcal{G} . I.e., $n,n' \in \mathcal{G}$ and $\forall a \in \mathcal{G}$

n*a = a*n = a and n*a = a*n' = a.

Consider a*n'= a. By Theorem A1.3.2, this ⇒

n' = a * a $= n \quad (neverse) . \square$

Theorem A1. 3.5. A given element of a group has only one reverse.

Proof. Let $a \in \mathcal{Y}$, where \mathcal{Y} is a group. Suppose \hat{a} and \hat{a} are both reverses of a. I.e., \hat{a} , $\hat{a} \in \mathcal{Y}$ and

a + a = a + a = n and a + a = a + a = n.

Consider axt = n. By Theorem A1.3.2, His ⇒

 $\overset{\star}{a} = \overset{\star}{a} + n$ $= \overset{\star}{a} \quad (\text{newhal}) \cdot \square$

Theorem A1, 3.6. Let n be the identity element of

Proof. For any element a, the reverse axiom requires $a \neq \hat{a} = n$. Taking a = n, we get $n \neq \hat{n} = n$. By Theorem A1.3.2, this \Rightarrow

 $n \in (neverse)$. \square

Theorem A1, 3, 7, Let a be any element of a group.

Proof. The reverse axiom requires $\hat{a}*a=n$. By Theorem A1.3.2, this \Rightarrow

 $a = \frac{r}{\delta * n}$ $=\overline{\hat{a}}$ (neutral). \Box

Another useful vesult is that the reverse of a combination is the combination of the veverses in the veverse order. Formally,

Theorem A13.8. Let a and b be any two elements of a group. $\frac{r}{a * b} = b * a'$

$$(a*b)*(b*a) = a*(b*(b*a))$$
 (associativity)
= $a*((b*b)*a)$ (associativity)
= $a*(n*a)$ (neverse)
= $a*a$ (rewind)
= n (neverse).

In the same way,

But these are the two properties that the reverse axiom requires. i. 6*6 is an inverse element to a*b. In view of Theorem A1,3,5, it is the reverse. \square

Since a group I is first of all a set, it is easy to consider subsets of I. An important greation is whether or not the subsets are themselves groups. In general, they are not. E.g., consider the subset I - {n}. The following terminology will be useful in this regard.

Definition A1.3.2. Let H be a nonempty [proper] subset of a group D. Then if H is a group w.r.t. the comb-ination function of S, it is called a [proper] subgroup of S.

Several comments are in order here. The bracket device is a standard way to make two similar statements at once. To get the first statement, leave out the the bracketed material. To get the second statement, include the bracketed material, Since HCY, HXHC YXY, and there is no problem in evaluating the combination function of Y at elements of H. However, the function values, which necessarily belong to Y, are not automatically in H; i.e., for an arbitrary subset, closure could fail.

theorem A1.3.9. Let I be a group with neutral element n. Then the singleton En's is a subgroup of I.

Proof, Exercise A1.3.4. [

Exercise A1.3.5. Refer to Theorem A1.3.9. Is {in} a proper subgroup of D?

The following theorem provides a useful criterion sufficient for a subset to be a subgroup.

Theorem A1.3.10, Let H be a nonempty subset of a group H, H a, $h \in H$, $A \times b \in H$,

then H is a subgroup of J.

Proof. As noted in the discussion following Definition A1.3.2, the contination function for I makes sause for H. Moreover, the associativity equivement is: automatically met. Assume that

(*) $a,b \in A \Rightarrow a*b \in A$,

Let $a \in H$ and choose b = a. Then $(*) \Rightarrow a * \tilde{a} \in H$. But $a * \tilde{a} = n$, the neutral element of H. Thus, $n \in H$.

Next, let $b \in \mathcal{H}$ and choose a = n. Then $(x) \Rightarrow n * b \in \mathcal{H}$. But n * b' = b'. Then $b \in \mathcal{H} \Rightarrow b' \in \mathcal{H}$, so \mathcal{H} contains the reverses of all of its elements. The reverses necessarily exist because of the properties of the "parent" group \mathcal{L} . Finally, let $a, b \in H$. Then by the preceding step, $b \in H$. Letting b' play the role of b in (*), we get $a*b \in H$.

But by Theorem A3.1.7, $\stackrel{r}{b} = b$. Thus, $a*b \in A$; and we have that A is closed under the combination operation, \Box

Supplementary Reading

BIRKHOFF and MAC LANE, A Survey of Modern Algebra

BOWEN and WANG, Introduction to Vectors and Tensors, Vol. 1, Linear and Multilinear Algebra

MICHEL and HERGET, Mathematical Foundations in Engineering and Science

MOSTOW, SAMPSON, and MEYER, Fundamental Structures of Algebra

NOLL, Finite - Dimensional Spaces

ODEN, Applied Functional Analysis

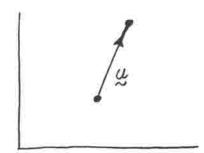
SIMMONS, Introduction to Topology and Modern Analysis

Chapter A2

Linear Spaces

A2.1. Definition of a Linear Space

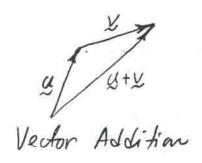
For the purpose of motivation, recall your experience with vectors in the <u>Euclidean plane</u>. Here a vector is just the Linected line segment or arrow that goes from one point in the plane to another point in the plane.



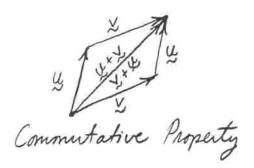
We will follow the common practice in mechanics of senoting vectors by bold face lower case letters.

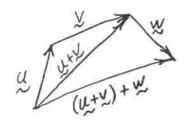
Two vectors in the plane are regarded as <u>equal</u> if they are parallel, point in the same direction, and have the same length. The actual location of a vector in the plane is not important, at least for most purposes.

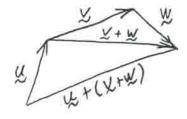
The addition of vectors is defined by the pullelogram rule on the "tip to tail" method depicted below. Note that the sum is itself a vector; i.e., vector addition is closed.



The addition of vectors is "immediately seen" to be commutative (u+v=v+u) and associative (u+v)+w=u+(v+w). The "proofs" of these properties are pictured next.







Associative Property

In this scheme, any two distinct points in the plane define a vector. What if the two points are the same? This means the arrow collapses to a single point. We naturally think of this peculiar vector as having

zero langth, but we cannot assign it a direction. It is called the zero vector and is denoted by Q. Because of its zero langth, it clearly has the property that for any vector ω , $Q+\omega=\omega+Q=\omega$. Thus, for vector addition, the zero vector plays the role of a neutral element.

Given any monzero vector u, we can generate a new vector by simply reversing the direction of u. This new vector is called the <u>negative</u> of u, and is denoted by -u.

u / -u

Clearly, the negative has the property that (-u)+u=u+(-u)=2. If we take the negative of Q to be Q itself, then we can say that all vectors have this property. Thus, for vector addition, the negative vectors play the roles of severse elements.

In comparison with Definition A1.3.1, we can say that the set of plane vectors is a group w.r.t. vector addition.

You also barned another operation with vectors, Given any nonzero vector u and any nonzero real number a,

Hen the <u>scalar multiple</u> of u by x, denoted by xu, is defined to be the vector which is parallel to u, points in the same or opposite Linetian as u according as x>0 or x<0, and has length |x||u|, where |u| is the length of u, If either x=0 or u=2, then x=0.

This operation also has important peoperties. First, we note the obvious notational delights that for any vector w, 1u = u and (-1)u = -u. It is important to realize that these are theorems — not notational agreements.

With a little more effort, we can see, as surely you have at some point in your studies, that for any \propto , $\beta \in \mathbb{R}$ and for any vectors y and χ , scalar multiplication is

associative, $(XB)u = \alpha(Bu)$,

distributive w.r.t. scalar addition, (x+B) w = xw + Bw,

and distributive w.r.t. vector addition, $\alpha(u+v) = \alpha u + \alpha v$.

In our study of continuum mechanics, we shall encounter many different sets and operations that have essentially the same structure that we have observed above for vectors in the Euclidean plane. Thus, it will be efficient to make the following generalization.

V, equipped with two special operations, The first operation, called addition, is a function in VXV to V, Senoted here

 $(\underline{u},\underline{v}) \mapsto \underline{u} + \underline{v}$

with the function value u+v called the sum of u and v, which satisfies the following axioms:

(A1) Associativity. Yu, V, W & V

(x+x)+x=x+(x+x);

(A2) Existence of a zero element, ∃ 0∈V >

Q+W=W+Q=W YWEV;

(A3) Existence of regative elements. For each $u \in V \exists$

 $-\omega + \omega = \omega + (-\omega) = \emptyset$

(A4) Commutativity. Yu, v EV

u+v=v+u.

The second operation, called scalar multiplication, is a function

vectors. Often, real vector space, and then the elements of V are called

on $\mathbb{R} \times \mathbb{V} \xrightarrow{to} \mathbb{V}$, denoted here by $(\alpha, u) \longmapsto \alpha u$

of U by &, which satisfies the following axioms:

(SMI) Associativity. Yx, B & R and Yu & V

(aB) = a (Bu);

(SM2) Distributivity w.r.t. scalar addition. ∀a, β ∈ IR and ∀ u ∈ V

 $(\alpha+\beta)u = \alpha u + \beta u$;

(SM 3) Distributivity w.r.t. addition in V. ∀ x ∈ R and

 $d\left(\underbrace{u}+\underbrace{v}\right)=d\underbrace{u}+d\underbrace{v};$

(SM4) YueV

1 = u.

More generally, R in the above Sefinition could be any "scalar field" IF. We would then speak of a linear space over the field IF. When IF is the complex number field I, then we have a complex linear space. We will use only real

linear spaces, and the adjective real will usually be amitted. Generally, the elements of linear spaces will be denoted by lower case bold face letters, and light fore letters will be used for elements of R, sometimes, entranched asstom or the likelihood of confusion will dictate otherwise,

If course, it is no accident that a linear space is defined in such a way that it is a commutative group wir.t. addition. Speaking of addition, note that the + on the S.h.s, of (SM2) refers to ordinary real number addition, while on the r.h.s. + refers to the addition operation on the linear space. If this causes confusion, return to GO and ask for a \$200 furtion refund.

We have already seen that the set of vectors in the Euclidean plane is a linear space. Another important and femilian example is given in the following theorem.

Theorem A2.1.1. Let n be a strictly positive integer. Consider the set of all list if length n whose elements are real numbers; i.e.,

 $\mathbb{R}^{\eta} := \left\{ (u_1, u_2, \dots, u_n) : u_i \in \mathbb{R}, i = 1, 2, \dots, n \right\}$

Write $\mathcal{U} = (u_1, u_2, \dots, u_n)$, $\mathcal{V} = (v_1, v_2, \dots, v_n)$, etc., and

[&]quot;I.h.s. stands for "left-hand side"; sunilarly n.h.s. means "right-hand side".

define addition, scalar multiplication, and equality by

 $\mathcal{L} + \mathcal{L} = (\mathcal{L}_1 + \mathcal{L}_1, \mathcal{L}_2 + \mathcal{L}_2, \dots, \mathcal{L}_n + \mathcal{L}_n),$ $\mathcal{L} \mathcal{L} = (\mathcal{L}_1, \mathcal{L}_2, \dots, \mathcal{L}_n),$

 $\mathcal{U} = \mathcal{V} \iff u_i = v_i, i = 1, 2, \dots, n$.

This system is a real linear space. It is called n-dimensional number space?

Proof. First, we examine the addition operation. Obviously, the Lomain is $\mathbb{R}^n \times \mathbb{R}^n$, the codomain is \mathbb{R}^n , and the evaluation rule is unambiguous. Thus, by Theorem A1.2.1, the operation of addition is, indeed, a function on $\mathbb{R}^n \times \mathbb{R}^n$ to \mathbb{R}^n .

Note that this is consistent with Definition A1.1.8.

n-dimensional authoretic space, n-dimensional condinate space,

Next we check the group axioms for addition. $(u+v)+w=((u_1+v_1)+w_1, \cdots, (u_n+v_n)+w_n)$ (definition of + in \mathbb{R}^n) = $(u_1+(v_1+w_1), --, u_n+(v_n+w_n))$ (associativity of + in IR and refinition of = in IRⁿ) = u + (v + w) (refinition of + in IRⁿ). (A2): Define $Q \in \mathbb{R}^n$ to be Q = (0,0,-,0). Then $\forall u \in \mathbb{R}^n$ $0 + u = (0 + u_1, \dots, 0 + u_n)$ (definition $g + u R^n$) $=(u_1,--,u_n)$ (property of $0 \in \mathbb{R}$ and definition of = in \mathbb{R}^n) = " (notation).

Similarly, $u+\varrho=u$.

(A3): Given
$$U_{+} = (U_{+}, U_{2}, \dots, U_{n}) \in \mathbb{R}^{n}$$
, define $-U_{+} \in \mathbb{R}^{n}$ to be $-U_{+} = (-U_{+}, -U_{2}, \dots, -U_{n})$. Then

$$-U_{+} + U_{+} = (-U_{+} + U_{+}, \dots, -U_{n} + U_{n}) \quad (\text{definition of } + \text{ in } \mathbb{R}^{n})$$

$$= (O_{5} - \dots, O) \quad (\text{property of negatives in } \mathbb{R} \text{ and } U_{+} + \text{ in } \mathbb{R}^{n})$$

$$= Q \quad (\text{definition of } Q \in \mathbb{R}^{n}).$$
Similarly, $U_{+} + (-U_{+}) = Q$.

$$(A4):$$

$$U_{+} = (U_{+} + V_{+}, \dots, V_{n} + V_{n}) \quad (\text{definition of } + \text{ in } \mathbb{R}^{n})$$

$$= (V_{+} + U_{+}, \dots, V_{n} + U_{n}) \quad (\text{commultivity of } + \text{ in } \mathbb{R} \text{ and } \text{ definition of } = \text{ in } \mathbb{R}^{n})$$

$$= V_{+} + U_{+} \quad (\text{definition of } + \text{ in } \mathbb{R}^{n})$$

Scalar multiplication can be checked out in a similar fashim, and this constitutes Exercise A2.1.1. [

Exercise A2.1.2. Consider the set of all real numbers IR. Denote them by x,y, etc., Define addition and scalar multiplication by ordinary real number addition and multiplication. Of course, the scalars are also elements of IR, to minimize confusion, Senote them by x, B, etc., Show that this system is a linear space. Is it different from IR' of Theorem A2.1.1?

Exercise AZ.1.3. Consider the set of all continuous, real-valued functions on the closed interval $[a,b] := \{x \in \mathbb{R} : a \le x \le b\}$; i.e.,

C([a,b]) = {f: f is a continuous function on [a,b] to IR}

Define addition, scalar multiplication, and equality by

(f+g)(x) = f(x) + g(x), $x \in [a,b]$,

 $(\alpha f)(x) = \alpha [f(x)], \quad x \in [a,b]$

and

 $f = g \iff f(x) = g(x) \ \forall x \in [a,b],$

Show that C°([a,b]) is a linear space. Feel free to how on theorems from calculus.

Since a linear space is a group w.r.t. addition (with 2 and -u playing the roles of the neutral and reverse elements, respectively), all of the results of \$A1.3 apply. From Theorems A1.3.4 and 5, we have

Theorem A2.1.2. A given linear space has only one zero element. A given element of a linear space has only one regative.

In the context of a linear space, the cancellation property (Theorem A1,3,3) becomes

14. Footnote 1 on p. A1. 2.13.

Theorem A2.1.3. For U, V, and w elements of a linear space $u+v=u+w \Rightarrow v=w$.

Since addition in a linear space is commutative, we so not bother to formally state the companion result that

 $\chi + \omega = \chi + \omega \implies \chi = \chi$

and, of course, this can be written in still other variants.

Theorems A1,3.6 and 7 give us the following interesting property of negatives

Theorem A2.1.4. Let u be any element of a linear space, and let 2 be the zero element. Then

-0 = 0 and --(-u) = u.

For stating the remaining consequences of the group structure, it will be convenient to introduce

Definition A2.1.2. The difference U-x of two elements u and x of a linear space is the sum of u and -x; i.e.,

 $\mathcal{U}-\mathcal{V} = \mathcal{U} + (-\mathcal{V}).$

The operation (U,V) - U-V is called subtraction.

If course, by closure the difference of two elements belows to the underlying linear space. By Axiom (A3), the difference has the following property.

Theorem A2.1.5. Let w be any element of a linear space. Then w-w=0.

Note that the above result is not true just because it looks right in terms of our experience with minus signs.

Theorem A1.3.8 has the following interpretation for linear spaces.

Theorem A2.1.6. Let a and & be any two elements of a linear space. Then

 $-(\underbrace{\omega}+\underline{\vee})=-\underbrace{\omega}+(-\underline{\vee})=-\underline{\omega}-\underline{\vee}.$

Theorem A1,3.2 gives us

Theorem AZ.1.7. Let V be a linear space. Then for any $u, v \in V$, \exists a unique $x \in V \ni$

 $\mathcal{L} + \mathcal{Z} = \mathcal{L}$.

In fact,

x = x + (-x) = x - x.

Finally, we note that in a linear space when equals are added to or subtracted from equals, the results are equal.

Theorem A2.1.8, Let S, t, u, and & be elements of a linear space. Then

 $S = \frac{t}{2} \text{ and } \mathcal{U} = \mathcal{V} \implies S + \mathcal{U} = \frac{t}{2} + \mathcal{V} \text{ and } S - \mathcal{U} = \frac{t}{2} - \mathcal{V}$

Proof, The result for sums is simply the linear space counterpart of Theorem A1.3.1.

To get the subtraction version, we note that u=x means that u and v are really the same element of the linear space; $u=v \Rightarrow -u=-v$, Then by the result for sums

S=t and $\mathscr{U}=\mathscr{V}$ \Rightarrow $S+(-\mathscr{U})=t+(-\mathscr{V})$

 $\Rightarrow 5 - \mu = \pm - 2$ (Definition A2.1.2).

The scalar multiplication analogue of the above result is

of a linear space. Then

 $\alpha = \beta \text{ and } \mathcal{U} = \mathcal{V} \Rightarrow \alpha \mathcal{U} = \beta \mathcal{V}.$

Proof. Exercise AZ.1, 4. [

(remaining)

For the properties of scalar multiplication, we have nothing to fall back on except for our experience with particular examples such as vectors in the Euclidean plane. The more common properties are gathered together in

Theorem A 2.1.10. Let u and v be any two elements of a real linear space, and let & and B be any two real numbers. Then

(iii)
$$\Delta u = Q \Rightarrow \underline{either} \Delta = 0 \underline{ou} u = Q$$
;

(iv)
$$-(\alpha \underline{u}) = \alpha(-\underline{u})$$
;

$$(v) - (\alpha \underline{\omega}) = (-\alpha)\underline{\omega};$$

$$(vi) - \omega = (-1) \omega ;$$

(Viii)
$$\alpha(\underline{u}-\underline{v}) = \alpha\underline{u} - \alpha\underline{v}$$
.

Proof. (i):
$$u = 1u$$
 (Axiom (SM4))

It is tempting to conclude at this point on the Basis of (A2) and the uniqueness of Q that Ou = Q. However, the Q in (AZ) must work & elements; i.e., (Do not have y= y+ou to conclude that Ou=Q) W= U+Q, X= V+Q, W= W+Q, etc,

But as far as we know now, Ou is only "a zero for u". To continue, we subtract a from both sides of the above result to get

= Ou (zero element).

For reasons which will become apparent, it is convenient to prove the remaining assertions in an order which differs

from that in which they were stated.

(V):
$$Q = \alpha U + [-(\alpha U)]$$
 (regative)

$$\Rightarrow (-\alpha)U + Q = (-\alpha)U + [-(\alpha U)] +$$

 $= \alpha \mathcal{L} + \left[- (\beta \mathcal{L}) \right] \quad (v)$

$$= \alpha(y - \beta y) \quad (Definition A2,1,2).$$
(iv): Take $\alpha = 0$ in (vii) to get
$$(0-\beta) y = 0y - \beta y \qquad (zero in R)$$

$$= 0y + [-(\beta y)] \quad (Definition A2,1,2)$$

$$= 0 + [-(\beta y)] \quad (i)$$

$$= -(\beta y) \quad (zero element).$$
(viii): $\alpha(y - y) = \alpha[y + (-y)] \quad (Definition A2,1,2)$

$$= \alpha y + \alpha(-y) \quad (distribution x)$$

$$= \alpha y + \alpha[(-1)y] \quad (vi)$$

$$= \alpha y + (-\alpha)y \quad (association y)$$

$$= \alpha y + (-\alpha)y \quad (property of R)$$

$$= \alpha y - \alpha y \quad (Definition A2,1,2).$$

$$\mathcal{A}(\mathcal{U} - \mathcal{U}) = \mathcal{A}\mathcal{U} - \mathcal{A}\mathcal{U}$$

$$\Rightarrow \alpha(0) = 0$$

(Theorem A 2.1.5).

(iii): Suppose that x u = Q, If $x \neq 0$, then we can use Theorem A2.1, 9 to multiply both sides by 1/x to get

$$\Rightarrow$$
 $l = 1/4 (x u)$ (ii)

$$= \%$$
 (SM4),

Thus, if $x \neq 0$, then $xy = Q \implies y = Q$.

Now assume that $\alpha w = 0$ and $w \neq 0$, We show that in this case $\alpha = 0$ by reductio and absundum; i.e., by argument by contradiction. Suppose $\alpha \neq 0$. Then we can use the argument above to conclude that w = 0 — contradicting the hypothesis that $w \neq 0$. $\therefore \alpha \neq 0$ is false, which we have that $\alpha = 0$.

Of course, if either $\alpha=0$ or $\alpha=0$ (or both) to begin with (and, in view of (i) and (ii), these are consistent with $\alpha\alpha=0$), then there is nothing to prove. \square

It is useful to realize that if we have a condidate for a linear space in the sense of having the elements and the operations of addition and scalar multiplication, then (i) and (vi) of Theorem A2.1.10 tell us exactly what the zero and the negatives must be; i.e., for any w

 $g = 0 \omega$ and $-\omega = (-1)\omega$.

In view of the uniqueness theorem, Theorem A2.1.2, there are no other possibilities. Often the candidate fails because the zero element on the negatives are not contained in the given set.

We generally will use the linear space oxioms and Theorems A2.1.3-10 without further comment. I.e., we will not meticulously write out every step such as we did in the proof of Theorem A2.1.10. However, one must, of course, always se able to so this,

Space V is said to be a linear subspace of V if it is closed w.r.t. addition and scalar multiplication; i.e.,

Often, simply subspace or linear menifold.

 $(51) \quad y, y \in \mathcal{S} \Rightarrow (y+y) \in \mathcal{S}$ $(52) \quad y \in \mathcal{S}, x \in \mathbb{R} \Rightarrow xy \in \mathcal{S}.$

Exercise A2.1.5. Refer to Definition A2.1.3. Prove that (S1) and (S2) together are equivalent to

 $u, v \in S$ and $\alpha, \beta \in R \Rightarrow (\alpha u + \beta v) \in S$

Theorem A2.1.11. A linear subspace of a linear space is itself a linear space.

Proof. Exercise AZ.1.6. [

a linear space, Then an element of V of the foun

a, d, + a, d, + -- + a k wk)

where $\alpha_1, \alpha_2, \dots, \alpha_k \in \mathbb{R}$, is said to be a linear combination of the subset $\{ \mathcal{U}_1, \mathcal{U}_2, \dots, \mathcal{U}_k \}$.

The set of all possible linear combinations of a given subset is of special interest.

of course, the operations of addition and scalar multiplication for the subspace are borrowed from the "parent" space.

Definition A2.1.5. Let & w, W, w, ..., Wk 3 C V, where V is a linear space. Then the set

 $Lsp\{u_1,u_2,...,u_k\}=\{u: u=\alpha_1u_1+\alpha_2u_2+...+\alpha_ku_k,\alpha_i\in\mathbb{R}\}$

is called the linear span of EU, Uz, -, Uh?

Theorem A2,1,12. Lsp & &, U, ..., uk } as defined above is a linear subspace of V.

Proof. Exercise AZ.1,7. []

More often, simply span, with the notation Sp {u, u, u, -, u, }.

Supplementary Reading

BISHOP and GOLDBERG, Tensor Analysis on Manifolds

BOWEN and WANG, Introduction to Vectors and Tensors, Vol.1, Linear and Multilinear Algebra

GEL'FAND, Lectures on Linear Algebra

GREUB, Linear Algebra

HALMOS, Finite-Dimensional Vector Spaces

LOOMIS and STERNBERG, Advanced Calculus

MARTIN and MIZEL, Introduction to Linear Algebra

MICHEL and HERGET, Mathematical Foundations in Engineering and Science

MOSTOW, SAMPSON, and MEYER, Fudamental Structures of Algebra

NAYLOR and SELL, Linear Operator Theory in Engineering and Science NICKERSON, SPENCER, and STEENROD, Advanced Calculus NOLL, Finite-Dimensional Spaces

NOMIZU, Fundamentals of Linear Algebra
ODEN, Applied Functional Analysis

A2.2. Finite-Dimensional Linear Spaces

There are many different ways to organize the material presented in this section, Consequently, some of our theorems will be definitions in alternative treatments and vice versa. We shall adopt what seems to us to be the most natural approach; it is not necessarily the most efficient,

Definition A2.2.1. A subset & W, Wz, -, Wk? of a linear space is said to be linearly dependent of at least one of the elements of the subset is a linear combination of the subset is a linear combination of the subset is linearly independent.

The next theorem provides a standard tool for proving that a subset is linearly dependent.

Theorem A2.2.1. A subset $\{u_1, u_2, \dots, u_k\}$ of a linear space is linearly dependent iff $\exists \alpha_1, \alpha_2, \dots, \alpha_k \in \mathbb{R}$ not all zero \ni

 $\alpha_1 \mathcal{U}_1 + \alpha_2 \mathcal{U}_2 + \cdots + \alpha_k \mathcal{U}_k = \mathcal{Q}$.

Proof. Suppose that { u, u, ..., uk? is linearly dependent.

Of course, the operations of addition and scalar multiplication from the underlying linear space are used to form the linear combination, of Definition A 2.1.4.

Y course, 2 is the zero element of the anderlying linear space.

Then are of the Q's, say Q_1 , can be written as a linear combination of the others; i.e., $\exists \beta_2, \beta_3, \cdots, \beta_R \in \mathbb{R} \ni$

U = B2 W2 + B3 W3 + --- + BE WE.

 $\Rightarrow \qquad \mathcal{U}_{1} - \beta_{2} \, \mathcal{U}_{2} - \beta_{3} \, \mathcal{U}_{3} - \cdots - \beta_{k} \, \mathcal{U}_{k} = \mathcal{Q} .$

 $X_1 U_1 + X_2 U_2 + \cdots + X_k U_k = Q$.

Conversely, suppose that I x's not all zero > this last equation holds. For definiteness, take x, \$0. Then the above equation >

$$\mathcal{U}_1 = -\frac{\chi_2}{\chi_1} \mathcal{U}_2 - \frac{\chi_3}{\chi_1} \mathcal{U}_3 - \cdots - \frac{\chi_R}{\chi_1} \mathcal{U}_R.$$

Here, one of the u's can be expressed as a linear contination of the others; and: the set $\{u_1, u_2, \cdots, u_k\}$ is linearly dependent, \square

The contrapositive of Theorem A2, 2, 1 provides a necessary and sufficient condition for linear independence. Eventhough it is less transparent than the prior content of this section, it is so useful that it is often taken as the starting point for this material.

Theorem A2.2.2. A subset {u, u, -, u, } of a linear space is is linearly independent iff the equation

 $\alpha_1 \overset{\vee}{\downarrow}_1 + \overset{\vee}{\lambda}_2 \overset{\vee}{\downarrow}_2 + \cdots + \overset{\vee}{\lambda}_k \overset{\vee}{\downarrow}_k = \overset{\circ}{Q} \quad (\alpha_i \in \mathbb{R})$

 $\Rightarrow \alpha_1 = \alpha_2 = \cdots = \alpha_k = 0.$

Proof. Suppose that

 $\chi_1 \, \mathcal{U}_1 + \chi_2 \, \mathcal{U}_2 + \cdots + \chi_k \, \mathcal{U}_k = \mathcal{Q} \implies \chi_1 = \chi_2 = \cdots = \chi_k = 0 \; .$

is linearly independent. Then

 $\alpha_1'''_1 + \alpha_2'''_2 + \cdots + \alpha_h \psi_k = 0 \implies \alpha_1 = \alpha_2 = \cdots = \alpha_h = 0.$

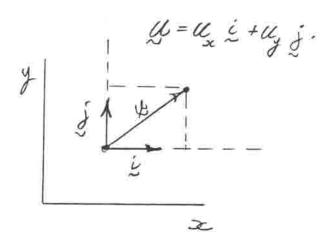
For if not, $\exists x's mot all zero <math>\ni x_1 \mathcal{U}_1 + \mathcal{U}_2 \mathcal{U}_2 + \cdots + \mathcal{U}_k \mathcal{U}_k = \mathcal{Q}$; and by Theorem A2.2.1, $\{\mathcal{U}_1, \mathcal{U}_2, \cdots, \mathcal{U}_k\}$ is linearly dependent — a cardiactichian. \Box

Exercise A2.2.1. Let $u \in V$, where V is a linear space. Show that if $u \neq Q$, then the singleton $\{u\}$ is linearly independent.

The following three easy results are often useful.

INSERT for p. A2.2.5

For much of the study of vectors in the Evelidean plane, it is advantageous to introduce unit vectors, say i and j, along distinguished orthogonal directions, say & and y, so as to express the typical vector u in the form



This, of course, is the first skep in the Levelopment of the subject of enalytic geometry. Our first skep in the generalization of this line of thought to abstract linear spaces is

Theorem AZ, 2, 3. Any finite subset (of a linear space) which contains the zero element is linearly dependent.

Proof, Exercise AZ. 2.2

Theorem A2,2,4. Any subset of a linearly independent subset is linearly independent.

Purf. Exercise A2.2.3. []

Theorem A2.2.5. Any finite subset (of a linear space) which contains a linearly dependent subset is linearly dependent.

Proof. Exercise A2, Z. 4. []

Definition A2.2.2. Let V be a linear space. A subset Egg, Ez, ..., En C V is said to be a basis for V if

{ €1, €2, ··· , €n} is linearly independent

and

(B2) { E1, E2, ..., En} spans 2 V in the sense that VCLSp { e, ez, ..., en }.

Our definition of linear dependence is restricted to finite sets, See the book by NAYLOR and SELL listed in the Supplementary reading for an extension of the notion to infinite sets. 2 G. Definition AZ. 1.5.

Necessarily, Lsp $\{e_1, e_2, \dots, e_n\} \subset V$, so we could have stated (B2) as $V = Lsp\{e_1, e_2, \dots, e_n\}$.

Theorem A2,2.6. All bases for a given linear space 1 contain the same number of elements.

Proof. Suppose that both $\{\xi_1, \xi_2, \dots, \xi_n\}$ and $\{f_1, f_2, \dots, f_m\}$ are bases for the linear space V, Since $\{\xi_1, \xi_2, \dots, \xi_n\}$ spans V and since $\{f_1, f_2, \dots, f_m\} \subset V$, \exists $\beta_{ij} \in \mathbb{R} \ni f_{ij} = \beta_{1i} \xi_1 + \beta_{2i} \xi_2 + \dots + \beta_{ni} \xi_n$, $i=1,2,\dots,m$.

Then for any choice of $X_1, X_2, \dots, X_m \in \mathbb{R}$, we have

$$\begin{aligned} \alpha_1 f_1 + \alpha_2 f_2 + \cdots + \alpha_m f_m &= \left(\sum_{j=1}^m \beta_{1j} \alpha_j \right) \underbrace{e}_1 + \left(\sum_{j=1}^m \beta_{2j} \alpha_j \right) \underbrace{e}_2 \\ &+ \cdots + \left(\sum_{j=1}^m \beta_{nj} \alpha_j \right) \underbrace{e}_n \end{aligned}$$

Consider the system of homogeneous linear algebraic equations $\sum_{j=1}^{m} \beta_{ij} \alpha_{j} = 0 , \quad i=1,2,\cdots,n,$

and think of the a's as unknowns. Suppose that m>n. Then there are more unknowns than equations, and it is always

It is possible that a given linear space will have no basis, this, in fact, is the common circumstance.

possible to find a montional solution. More precisely, $\exists \alpha_1, \alpha_2, \cdots, \alpha_m \in \mathbb{R}$ not all zero \exists

 $\sum_{j=1}^{m} \beta_{ij} x_{j} = 0 , i=1,2,\dots,n.$

With such a choice of the x's

 $\alpha_1 f_1 + \alpha_2 f_2 + \cdots + \alpha_m f_m = Q$,

which by Theorem A2,2,2 contradicts the linear independence of $\{f_1, f_2, \dots, f_m\}$. if $m \le n$.

Next suppose that m < n and reverse the roles of $\{ \subseteq_1, \subseteq_2, \cdots, \subseteq_n \}$ and $\{ f_1, f_2, \cdots, f_m \}$ in the above regument to conclude that $m \ge n$. The details are left as Exercise A2,2.5. Hence, m = n. \square

In view of Theorem AZ.Z. 6, the following definition is maningful-

Definition A2.2.3. A linear space V is n-dimensional if it contains a basis with nelements; If no such integer nexists, then the linear space is infinite-dimensional.

Of course, n is called the dimension of V, and we write n=:dim V.

See almost any text that treats linear algebraic equations. A particularly good one is HOHN'S Elementary Matrix Algebra.

² Often, simply finite-dimensional.

Symbols such as In are often used to senote an n-dimensimal space.

Exercise A2.2.6, Consider n-chinensimal number space Rn. Prove that the lists

en = (0,0, - -, 0,1)

constitute a basis for Rn. It then follows from Definition A2.2.3 that Rn is , indeed, n-chinensimal as its name suggests. The above basis is called the standard basis for Rn.

Exercise A2.2.7. In the spirit of Exercise A2.1.2, show that the singleton {1} is a linearly independent subset which spans R. Thus, R, when viewed as a linear space, is 1-dimensional.

The following result will not be a surprise.

Theorem A2,2,7. Let V_n be an n-dinansimal linear space. Then any subset of V_n which contains more than n elements is linearly dependent,

Proof. Consider $\{u_1, u_2, \dots, u_m\} \subset V_n$ with m > n. Since V_n is n-dimensional, it has a basis $\{e_1, e_2, \dots, e_n\}$. Since

 $\{ \mathcal{U}_1, \mathcal{U}_2, \cdots, \mathcal{U}_m \} \subset \mathcal{V}_n \subset Lsp \{ \mathcal{E}_1, \mathcal{E}_2, \cdots, \mathcal{E}_n \} , \exists \beta_{ij} \in \mathbb{R} \ni$ $\mathcal{U}_i = \beta_{1i} \mathcal{E}_1 + \beta_{2i} \mathcal{E}_2 + \cdots + \beta_{ni} \mathcal{E}_n , i=1,2,\cdots, m.$

We leave it as Exercise A2.2.8 for the student to continue the argument as in the proof of Theorem A2.2.6 to the conclusion that {U, U2, -- , Um} is linearly dependent. \square

The following theorem is helpful in finding a casis when the timension of the space is already known,

Theorem A2.2.8. Let V_n be an n-dimensional linear space. Then a subset $\{\xi_1, \xi_2, \cdots, \xi_n\} \subset V_n$ is a basis for V_n if it is linearly independent,

Proof, Suppose that { e, e, ..., en is a basis for Vn. Then it is linearly independent by Definition A2,2,2.

Conversely, suppose that $\{e_1, e_2, \dots, e_n\}$ is linearly independent, In order to prove that it is a basis for V_n , we must show that it speems V_n . Accordingly, let w be an arbitrary element of V_n and consider the subset $\{u, e_1, e_2, \dots, e_n\}$. By theorem A2.2.7, this subset of 11+1 elements must be linearly dependent, if by Theorem A2.2.1, $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$ not all zero $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$ not all zero $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$ not all zero $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$ not all zero $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$ not all zero $\exists x_0, x_1, x_2, \dots, x_n \in \mathbb{R}$

 $x_0 = + x_1 = 1 + x_2 = 2 + \cdots + x_n = 0$.

Now $x_0 \neq 0$. For if $x_0 = 0$, then the above equation would become

$$\alpha_1 \underset{\sim}{\varepsilon}_1 + x_2 \underset{\sim}{\varepsilon}_2 + \cdots + \alpha_n \underset{\sim}{\varepsilon}_n = 0$$

with not all of $\alpha_1, \alpha_2, \dots, \alpha_n$ zero; and by Theorem A2.2.2, this contradicts the linear independence of $\{\xi_1, \xi_2, \dots, \xi_n\}$. Since $\alpha_0 \neq 0$, we can manipulate

into $u = -\frac{x_1}{d_0} \underbrace{e}_1 - \frac{x_2}{d_0} \underbrace{e}_2 - \dots - \frac{x_n}{d_0} \underbrace{e}_n \in Lsp \underbrace{\xi e}_1, \underbrace{e}_2, \dots, \underbrace{e}_n \underbrace{\xi}_n.$

: 1/n C Lsp {e, ez, -, en}. 0

We shall have occasion to use the next two theorems.

Theorem A2.2.9. Let 11, be an 11-dimensional linear space.

Iku any linearly independent subset $\{e_1, e_2, \dots, e_m\}$ $\in V_n$ with m < n can be extended to a basis for V_n , i.e., \exists n-m elements e_{m+1} , e_{m+2} , \cdots , $e_n \in V_n \ni$ $\{e_1, e_2, \dots, e_m\}$ $\{e_m\}$, $\{e_m\}$,

Proof. Choose $\mathcal{L}_{m+1} \in \mathcal{V}_n \ni \mathcal{L}_{m+1} \notin Lsp \mathcal{L}_{2}, \mathcal{L}_{2}, \dots, \mathcal{L}_{m}$. This is always possible. For if not, then the linearly independent subset $\mathcal{L}_{1}, \mathcal{L}_{2}, \dots, \mathcal{L}_{m}$ spans \mathcal{V}_{n} ; consequently, if is a basis for \mathcal{V}_{n} , and n=m-a contradiction. The subset $\mathcal{L}_{2}, \mathcal{L}_{2}, \dots, \mathcal{L}_{m}$ is linearly independent. For

if not, it is linearly dependent; then since $\{\mathcal{L}_1, \mathcal{L}_2, \cdots, \mathcal{L}_m\}$ is linearly independent, it is easy to show as in the proof of Theorem A2.2.8 (The details are left as Exercise A2.2.9.) that $\mathcal{L}_{m+1} \in Lsp \{\mathcal{L}_1, \mathcal{L}_2, \cdots, \mathcal{L}_m\}$ — a contradiction. If m+1=n, then by Theorem A2.2.8, we are done.

If m+1 < n, then we repeat the procedure. Choose $\mathbb{E}_{m+2} \in V_n \ni \mathbb{E}_{m+2} \notin Lsp \{\mathbb{E}_1, \mathbb{E}_2, \cdots, \mathbb{E}_m, \mathbb{E}_{m+1}\}$. Then $\{\mathbb{E}_1, \mathbb{E}_2, \cdots, \mathbb{E}_m, \mathbb{E}_{m+1}, \mathbb{E}_{m+2}\}$ is linearly integrandant. If m+2=n, we are done.

leady, after a total of precisely n-m such steps, we have an extended fasis. (Exercise A 2.2.10. Devise a matematical induction argument to show that the procedure never breaks down.) []

A useful complement to theorem A2.2.9 is given by

Theorem A2,2.10. Let $11 \pm \{Q\}$ be a linear subject of an 11-dimensional linear space V_n . Then 11 is 1n-dimensional with 11 ± 11 , and \exists a bosis $\{E_1, E_2, \cdots, E_m, E_{m+1}, E_{m+2}, \cdots, E_n\}$ for $V_n \ni \{e_1, e_2, \cdots, e_m\}$ is a bosis for 20.

Proof. Since U is a linear subspace, it is not empty; and since U + {Q}, U contains a nonzerol, say E, . By Exercise A2.2.1, {E, } is linearly independent. i. if U < Lsp {E, }, {E, } is a basis for U; and U is 1-dimensional. By Theorem A2.2.9, {E, } and be enlarged to a basis for Un; and the proof is complete.

If U \$ Lsp {e, },] = e = U = e & Lsp {e, }. It is

lasy to show (The Setails are left as Exercise A2.2.11.) that $\{\mathcal{L}_1,\mathcal{L}_2\}$ is linearly independent. If $\mathcal{U} \subset L_{SP}\{\mathcal{L}_1,\mathcal{L}_2\}$, $\{\mathcal{L}_1,\mathcal{L}_2\}$ is a basis for \mathcal{U} ; and \mathcal{U} is 2-dimensional. Again $\{\mathcal{L}_1,\mathcal{L}_2\}$ can be extended to a basis for \mathcal{V}_n , and we are some.

If $U \neq Lsp \{e_1, e_2\}$, then we repeat the proceduce. Since $U \subset V_n$, which is n-dimensional, there can be no more than n such steps. \square

By Definition A2.2.2, lack element of a finite-dimensional clinear space can be expressed as a linear combination of basis elements. The following terminology will prove useful when we take advantage of this fact.

Definition AZ. Z. 4. Let Ee, e, e, en! le a Basis for an II-dinvensional linear space Vn. Then the scalars!

U, u, u, ..., un in the representation

 $u = u^{\prime}e_1 + u^{\prime}e_2 + \cdots + u^{\prime\prime}e_n$

of $u \in V_n$ are the components of u w.r.t. the basis $\{e_1, e_2, \cdots, e_n\}$.

Here, scalar is just a synonym for real number. More generally, scalars are the elements of the scalar field which the linear space is over.

vector components. Sometimes, coordinates;

We used superscripts nother than subscripts on the components for a reason which will become clear later. At this point, it seems like a silly practice; adopt it anyway.

The concept of components would be useless if it were not for

Theorem A2.2.11. The components of a given element of a finite-dimensional linear space w.r.t. a given basis are unique.

Proof. Let $\{e_1, e_2, \cdots, e_n\}$ be a basis for an n-dimensional linear space V_n . Let $u \in V_n$, and suppose that we have both

 $u = u^{\prime}\underline{e}_{1} + u^{2}\underline{e}_{2} + \cdots + u^{\prime}\underline{e}_{n}$ and $u = \overset{\star}{u}^{\prime}\underline{e}_{1} + \overset{\star}{u}^{2}\underline{e}_{2} + \cdots + \overset{\star}{u}^{n}\underline{e}_{n}$.

These two equations >

 $(\ddot{a}'-u')e_{n} + (\ddot{a}'-u')e_{n} + \cdots + (\ddot{a}''-u'')e_{n} = 0$

which in turn >

 $\ddot{a}'=a', \quad \ddot{a}^2=a^2, \quad \cdots, \quad \ddot{a}^n=a^n,$

by Theorem A2.2.2 since {€1, €2, -.., €n} is linearly independent. □

The next result shows that addition, scalar multiplication, and equality of elements of a finite-dimensional linear space

correspond to addition, scalar multiplication, and quality of their components wirit, a basis,

Theorem A2.2,12, Let $\{e_1, e_2, \dots, e_n\}$ be a basis for an n-dimensional linear space V_n . Let $u, x \in V_n$ so that

u= u'e, + u'e, + u'en , x=v'e, + v'e, + v'en. Then

(i) For $\alpha, \beta \in \mathbb{R}$,

 $\alpha \mathcal{U} + \beta \mathcal{V} = (\alpha u^1 + \beta v^1) \underline{e}_1 + (\alpha u^2 + \beta v^2) \underline{e}_2 + \cdots + (\alpha u^n + \beta v^n) \underline{e}_n$

(ii) Q = 0e, +0e, + -- +0en;

(iii) $u=v \iff u^i=v^i, i=1,2,\cdots,n$.

Proof. Exercises A2.2.12-14. [

This easy theorem is actually quite profound. In general, two "matternatical structures" are "essentially identical if they can be put into one-to-one correspondence in a way that preserves their structure, i.e., they have the same mathematical structure even though the interpretations of their elements may differ. Technically, it is said that two such structures are "isomorphic", and the one-to-one correspondence which connects them is called an "isomorphism". In the case of a linear space, the (algebraic) I.e., in transparent notation, (xu+Bx) = xu+Bvi.

structure is provided by the operations of addition and scalar multiplication, and the above remarks lead us to

Definition A2.2.5. Two real linear spaces U and V are isomorphic if \exists a one-to-one and onto function $f: \mathcal{U} \rightarrow \mathcal{V}$ which is additive,

 $f(\alpha + x) = f(\alpha) + f(x) \quad \forall \alpha, x \in \mathcal{U},$

and homogeneous,

 $f(\Delta u) = \chi f(u) \quad \forall u \in \mathcal{U} \text{ and } \forall x \in \mathbb{R}$.

Any such function f is an isomorphism of U onto V.

Of course, He additivity is the presentation of the operation of addition, while the homogenesty is the preservation of scalar multiplication.

Theorem A2.2.13, Every n-dimensional linear space Vn is isomorphic to n-dimensional number space Rn.

Proof. Let $\{\xi_1, \xi_2, \dots, \xi_n\}$ be a basis for V_n , Then for $\xi \in V_n \exists u', u^2, \dots, u^n \in \mathbb{R} \ni$

 $U = u' \underline{e}_1 + u' \underline{e}_2 + \cdots + u'' \underline{e}_n$

With this representation understood, define $f: V_n \to \mathbb{R}^n$ by f(u) = (u', u', --, u'').

Since for given it the viare unique by Theorem AZ.Z.11, f is, indeed, a function, of. Theorem A1. 2.1.

To see that f is one-to-one, let $\mathring{\mathcal{C}}, u \in V_n$ and suppose that $f(\mathring{u}) = f(u)$. This \Rightarrow

 $(\ddot{a}', \ddot{a}^2, \dots, \ddot{a}^n) = (u', u^2, \dots, u^n)$

 $\Rightarrow \dot{u}^i = u^i, i=1,2,...,n$ (equalify in \mathbb{R}^n)

⇒ # = # (Theorem A2,2.12 (iii)).

i. by Definition A1,2.3, f is one-to-one.

To see that f is onto, let $(u', u', --, u'') \in \mathbb{R}^n$. Then $u := u'e_1 + u^2e_2 + -- + u''e_n \in Lsp \S e_1, e_2, --, e_n \rbrace = V_n$, and by the definition of f, $f(u) = (u', u^2, --, u'')$. i. by Definition A1, 2.2, $\mathbb{R}^n \subset f(V_n)$. But by the definition of f, $f(V_n) \subset \mathbb{R}^n$; so $f(V_n) = \mathbb{R}^n$, and f is onto.

By Theorem A2.2.12 (i), for $x, \beta \in \mathbb{R}$ and $u, v \in V_n$, $\angle x + \beta v = (\alpha u^1 + \beta v^1) \underbrace{e}_1 + (\alpha u^2 + \beta v^2) \underbrace{e}_2 + \cdots + (\alpha u^n + \beta v^n) \underbrace{e}_n.$

Then by the definition of f, $f(\alpha u + \beta v) = (\alpha u^1 + \beta v^1, \alpha u^2 + \beta v^2, \dots, \alpha u^n + \beta v^n)$ $= (\alpha u^1, \dots, \alpha u^n) + (\beta v^1, \dots, \beta v^n) \quad (\text{addition in } \mathbb{R}^n)$ $= \alpha (u^1, \dots, u^n) + \beta (v^1, \dots, v^n) \quad (\text{scalar multiplication in } \mathbb{R}^n)$ $= \alpha f(u) + \beta f(v) \quad (\text{definition of } f).$

This gives both additivity ($\alpha = \beta = 1$) and homogeneity ($\beta = 0$).

Hence, f is an isomorphism and Vn and R" are isomorphic. [

Note that in the above proof each different choice of basis for Vn (and there are infinitely many) sends to a different isomorphism of between Vn and Rn. This messy feature is one of the reasons that we do not just restrict our study of n-dimensional linear spaces to Rn from this point forward.

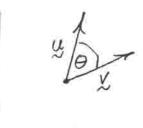
Supplementary Reading Same as for \$ AZ.1

A2.3. Inner Product Spaces. Normed Linear Spaces

Recall again your experience with the <u>Euclidean plane</u>. A great help to the study of the geometry of the plane is the concept of the <u>inner product</u> of two <u>vectors</u> in the plane, which is defined as

w.x = |x/1/x/cos0,

where |4| and |X| are the lengths of the vectors, and to is the (smaller) angle between them.



Here, we have an operation which takes two vectors and produces a scalar. What are its properties?

Obviously, it is commutative; i.e.,

~ ~ ~ = × · ~

It is also homogeneous in the sense that for $\alpha \in \mathbb{R}$ $(\alpha u) \cdot v = \alpha (u \cdot v)$.

(have)
To see this, we three cases to consider: 0.0, 0.0, 0.0.

Let us look at the most difficult case, which is 0.0.

In this case 0.0 points in the direction opposite to 0.0.

By Sefinition,

$$(\alpha u) \cdot v = |\alpha u||v|\cos(\pi - \theta)$$

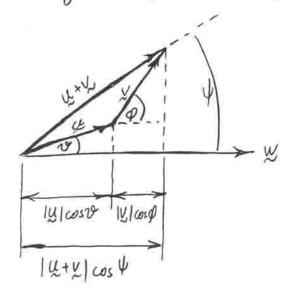
=
$$(-|\alpha|)||\alpha|||\alpha||\cos\theta$$
 (associativity and commutativity in \mathbb{R})
= $|\alpha||\alpha||\alpha||\alpha|$ (definition of $|\alpha|$ for $|\alpha|$)

$$= \alpha(\underline{u}.\underline{v})$$
 (definition of $\underline{u}.\underline{v}$).

The inner product of plane vectors is distributive w.r.t.

 $(\alpha + \gamma) \cdot \omega = \alpha \cdot \omega + \gamma \cdot \omega$

A "proof" of this fact is depicted below.



Finally, since $\psi \cdot \psi = |\psi|^2$, we see that the inner product is positive refinite in the sense that

 $\underline{\mathcal{U}} \cdot \underline{\mathcal{U}} \geq 0$ and $\underline{\mathcal{U}} \cdot \underline{\mathcal{U}} = 0$ iff $\underline{\mathcal{U}} = \underline{\mathcal{Q}}$,

Those are the key properties of the inner product of vectors in the Euclidean plane. For the general case, we have

Definition A2.3.1. A real inner product = pace is a val linear = pace, say V, equipped with a function on VXV to R, denoted by

 $(\underline{u},\underline{v}) \mapsto \underline{u} \cdot \underline{v}$

¹ Sanetimes, Euclidean space.

² Mathematicians usually write (u, v) or <u, x> for u.x.

which satisfies the following axioms:

(II) Commutativity. ∀ u, x ∈ V

K·X = X· ₩ ;

(I2) Homogeneity. $\forall \alpha \in \mathbb{R} \text{ and } \forall \underline{u}, \underline{v} \in V$ $(\alpha \underline{u}) \cdot \underline{v} = \alpha (\underline{u} \cdot \underline{v});$

(I3) Distributivity w.r.t. addition. $\forall u, v, w \in V$ $(u+v) \cdot w = u \cdot w + v \cdot w$;

(I4) Positive Definiteness, Yuev

 $\psi \cdot \psi \ge 0$ with $\psi \cdot \psi = 0$ only if $\psi = Q$.

For a complex linear space, $u \cdot v \in \mathbb{C}$, and the commutativity axim is modified to $u \cdot v = \overline{x \cdot u}$. This makes $u \cdot u \in \mathbb{R}$, so the inequality in the positive definiteness axiom makes sense. Complex inner product spaces are often called limitary linear spaces.

We have already seen that the set of vectors in the

Often, scalar product or dot product.

Euclidean plane is an inner product space. Another important example is \mathbb{R}^n .

Theorem A2,3,1, The linear space

 $\mathbb{R}^{n} = \{ (u_{1}, u_{2}, \dots, u_{n}) : u_{i} \in \mathbb{R}, i = 1, 2, \dots, n \}$

with u.x defined by

 $\mathcal{U} \cdot \mathcal{V} = \mathcal{U}_1 \mathcal{V}_1 + \mathcal{U}_2 \mathcal{V}_2 + \cdots + \mathcal{U}_n \mathcal{V}_n$

the standard inner product.

Proof. It is don 2 that

 $(\underline{U},\underline{V}) \longrightarrow \underline{U}_1 \underline{V}_1 + \underline{U}_2 \underline{V}_2 + \cdots + \underline{U}_n \underline{V}_n =: \underline{U} \cdot \underline{V}$

defines a function on $\mathbb{R}^n \times \mathbb{R}^n$ to \mathbb{R} . It is left for the student as Exercise A2.3.1 to show that axioms (I1)-(I3) are satisfied.

To verify the positive definiteness, we note that $u \cdot w = u_1^2 + u_2^2 + \cdots + u_n^2$

¹ G. Theorem AZ, 1, 1,

² G. Theorem A1.2.1.

Since there is no possibility of cancellation,

$$\mathcal{U}\cdot\mathcal{U}=\mathcal{U}_1^2+\mathcal{U}_2^2+\cdots+\mathcal{U}_n^2=0 \Rightarrow \mathcal{U}_{\epsilon}=0, i=1,2,\cdots,n \Rightarrow \mathcal{U}=0,$$

because the zero element of R" is (0,0, --, 0). Thus,

$$\underline{\mathcal{U}} \cdot \underline{\mathcal{U}} = 0$$
 only if $\underline{\mathcal{U}} = \underline{\mathcal{Q}}$,

and (I4) is met. [

The next exercise emphasizes that the standard inner product for Rn is not the only choice of inner product for Rn.

Exercise A2.3.2. Let the nxn real matrix [aij] be symmetric and positive definite. Look up the meanings of these terms in the matrix context if you do not know them. With reference to the notation of Theorem A2.3.1, show that

$$u \cdot v = \sum_{i=1}^{n} \sum_{j=1}^{n} \alpha_{ij} u_{i} v_{j}$$

defines an inner product for Rn.

Exercise A2.3.3. In the spinit of Exercises A2.1.2 and A2.2.7,

view R as a 1-dimensional linear space. Show that

$$\langle x, y \rangle = xy$$

defines an inner product on R. Here, the notation x y is avoided because of possible confusion with the product xy.

Since every M-dimensional linear space is isomorphic to \mathbb{R}^n (of Theorem A2, 2.13) and since \mathbb{R}^n can be made an inner product = pace in many ways, the following result is not unexpected.

Theorem A2,3,2. Every n-dimensimal linear space on can be assigned an inner product.

Proof. We make the same identification between V_n and \mathbb{R}^n that was used to establish that they are isomorphic. Let $\{\mathcal{L}_1,\mathcal{L}_2,\cdots,\mathcal{L}_n\}$ be a basis for V_n . Then for $u\in V_n$ \exists $u',u^2,\cdots,u^n\in\mathbb{R}$ \ni

 $u = u'e, + u^2e_2 + \cdots + u^n e_n.$

With this representation understood, Lefine

 $U_{\sim \sim}^{1} = u^{1}v^{1} + u^{2}v^{2} + \dots + u^{n}v^{n}.$

The remaining steps are identical to those used to prove

Theorem A2.3, 1, Here, the wearings of u+v and xu in terms of components follow from Theorem A2.2.12 as does the equivalence $u=0 \Leftrightarrow u^i=0$, $i=1,2,\cdots,n$.

If course, in the above proof, each different choice of basis can be expected to lead to a different inner product. Surprisingly, we shall see in \$ A2.5 that for a fairly large class of bases these inner products will be identical.

So far, all of our examples of inne product spaces have been finite-dimensional. The next exercise shows that infinite - dimensional spaces also can have an inner product.

Exercise A2.3.4. Consider the infinite-chinensianal linear space C°([a,b]) of Exercise A2.1.3, Show that

$$\langle f,g \rangle = \int_{a}^{b} f(x) g(x) dx$$

defines an inner product for C°([a,b]). (Strictly speaking, we haven't shown that this space is infinite-chinensional, but Lord spend too much time looking for a casis.)

The next three easy theorems provide useful usults that apply to all inner product spaces,

Theorem AZ.3,3, Let V be an inner product space. Their

$$\varrho \cdot \underline{\omega} = 0 \quad \forall \, \underline{\omega} \in \mathcal{V}.$$

In particular,

0.0 = 0.

Proof. Exercise A2.3.5. (Hint: Make intelligent choices of the variables in either the homogeneity property or the Listriorativity property.) [

Since we always have 0.0 = 0 in an inner product space, the positive definiteness axiom is sometimes stated in the Whendant fashion:

 $\forall \, \underline{u} \in V, \, \underline{u} \cdot \underline{u} \geq 0 \text{ with } \underline{u} \cdot \underline{u} = 0 \text{ if } \underline{u} = \underline{Q}.$

Theorem 12.3.4. Let $u \in V$, where v is an inner product space. Then

 $\tilde{\mathbf{x}} \cdot \tilde{\mathbf{x}} = 0 \quad \forall \, \tilde{\mathbf{x}} \in \mathcal{A} \Rightarrow \tilde{\mathbf{x}} = \tilde{\mathbf{0}}.$

Proof. Since $u \cdot v = 0$ must hold $\forall v \in V$, it must hold in particular for v = u. Hence, $u \cdot u = 0 \Rightarrow u = 0$ by the positive definiteness of the inner product. \Box

INSERT material on p.AZ.3.9a

In our motivational example of vectors in the Euclidean plane, $\psi = |\psi|^2$, where $|\psi|$ is the length of the vector ψ . This leads us to

Definition A2.3.2. The magnitude of an element us of

Often, length, modulus, norm.

INSERT for p. A2.3.9

As expected, of Theorems A1,3,1, A2,1.8, and A2.1.9, when equal elements of an inner product space are "dotted" with equals, the resulting scalars are equal.

an inner product space, Then

W= y and W= x - W· W = V·x.

Proof. Exercise A 2,3,6, []

Return to p. A2.3.9

of an inner product space is the is the nonnegative number

$$|\underline{u}| = \sqrt{\underline{u} \cdot \underline{u}} .)^{1,2}$$

Of course, this definition would not make sense if it were not for the fact that $u\cdot u \ge 0$ by the positive definiteness of the inner product.

Exercise A2.3.7. In the spirit of Exercises A2.1.2, A2.2.7, and A2.3.3, view R as the 1-dimensional inner product < pare with the inner product < x,y> = xy. Show that the magnitude of x is the absolute value of x. I.e., |x| = |x|; where on the 1.h.s. $|\cdot|$ means magnitude, while on the 1.h.s. $|\cdot|$ means absolute value. What good fortune!

It is important to realize that in the lifinition of magnitude and in most of what follows the underlying inner product space need not be finite-dimensional.

of course, the operation of calculating the magnitude of an element of an inner product space can be thought of as a function on the space to the nounegative reals. This is our viewpoint in stating the following important properties.

We will always use to to denote the positive square root.

Mothematicians usually write $\|\Psi\|$ for $|\Psi|$ (cf. Definition A2.3.3 and Theorem A2.3.9); engineers often use $|\Psi| = \Psi$.

Theorem A2.3.6. On any inner product space V, magnitude is homogeneous, positive

and positive definite,

(ii) | | | ≥0 + u ∈ V with | | | =0 iff u = Q.

Proof. Exercises A2,3,8 and 9. [

The following inequality has major applications in widely different contexts.

Theorem AZ. 3.7. (Inner Product - Mognitude Inequality) Let V be an inner product space. Then

| \(\cdot \cdot \) \(\lambda \) \(\cdot \) \(\cdo \) \(\cdot \) \(\cdot \) \(\cdot \) \(\cdot \) \(\cdo \) \(\cdot \) \(\cdot

Proof. If either W=Q or Y=Q, it follows from our provious results that the inequality above is satisfied as an equality. ... we assume that $W\neq Q$ and $Y\neq Q$ in the remainder of the proof.

By the positive definiteness of the inner product,

Almost always labled with some combination of the names CAUCHY, BUNYAKOVSKII, and SCHWARZ.

 $(\alpha \underline{u} + \beta \underline{v}) \cdot (\alpha \underline{u} + \beta \underline{v}) \ge 0 \quad \forall \alpha, \beta \in \mathbb{R}$

Expanding this (using distributivity, etc.), we find that

 $||\chi|^2 + 2\alpha\beta ||\chi| \cdot |\chi| + ||\beta|^2 ||\chi||^2 \ge 0,$

Choosing $x = |\chi|^2$ and $\beta = -u \cdot v$ yields $\left[|u|^2 |\chi|^2 - (u \cdot v)^2 \right] |\chi|^2 \ge 0.$

Since |v|2 >0, this >

 $|\underline{\omega}|^2 |\underline{v}|^2 \geq (\underline{\omega} \cdot \underline{v})^2$;

and since the square root function is monotonic, this >

 $|\underline{u}||\underline{v}| \ge \sqrt{(\underline{u}\cdot\underline{v})^2} = |\underline{u}\cdot\underline{v}| . \square$

It is interesting to note that when the inner product - magnitude in equality is applied to Rⁿ with the standard inner product (cf. Theorem A2.3.1) and to C^o[(a,b)] with the inner product introduced in Exercise A2.3.4, the following seemingly unrelated results are obtained:

$$\left| \sum_{i=1}^{n} a_i b_i \right| \leq \left(\sum_{i=1}^{n} a_i^z \right)^{1/2} \left(\sum_{i=1}^{n} b_i^z \right)^{1/2}$$

Here, of course, (.) 1/2:= V.

for any lists (a, az, ..., an) and (b, bz, ..., bn) of wal numbers; and

$$\left| \int_{a}^{b} f(x)g(x)dx \right| \leq \left(\int_{a}^{b} [f(x)]^{2} dx \right)^{k_{2}} \left(\int_{a}^{b} [g(x)]^{2} dx \right)^{k_{2}}$$

for any real-valued functions fand a continuous on [a,b]. This is indicative of the power of abstract methods.

In using the inner product - magnitude inequality, it is often helpful to ucoquize that

u·v ≤ | u·v | ≤ | u| | y |

in fact, we shall do this in our proof of

Theorem A2.3, B. (Triangle Inequality) Let V le on inner product space. Then

 $|\underline{u}+\underline{v}| \leq |\underline{u}|+|\underline{v}| \quad \forall \, \underline{u}, \underline{v} \in \mathcal{V}.$

Proof. $|u+v|^2 = (u+v) \cdot (u+v)$ (definition of 1.1 in v) $= |u|^2 + 2u \cdot v + |v|^2 \quad (properties of inner product)$

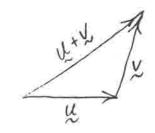
≤ |U|2+2|U·V|+ |V|2 (property of 1.1 in 1R)

\[
 \leq |\frac{1}{2} + 2 |\frac{1}{2} |\frac{1}{2}| + |\frac{1}{2}|^2 \quad \text{(inner product - Wagnitude in eq.)}
 \]

$$= (|\underline{u}| + |\underline{v}|)^{2} \quad (arithmetic)$$

$$\Rightarrow |\underline{u} + \underline{v}| \leq |\underline{u}| + |\underline{v}| \quad (anotonicity of \sqrt{\cdot}). \square$$

On recalling the rule for adding vectors in the Euclidean plane,



we see why the triangle inequality is so nemed. That it is instrinctively known to all school-aged drildren is a fact readily attested to by any residential homeowner with a fenceless corner lot.

Now that we have the triangle inequality, we are at a good place to introduce

Definition A2.3.3. A normed linear space is a real linear space, say V, equipped with a function on V to R+), senoted by $u \mapsto \|u\|$

Not necessarily an inner product space.

² \mathbb{R}^+ denotes the nonnegative reals; i.e., $\mathbb{R}^+:=\{x\in\mathbb{R}: x\ge 0\}$.

with the function value || W || called the norm of w, which satisfies the following axioms:

(NI) Homogeneity. YXER and YUEV

(N2) Positive Definiteness. VuEV

|| \(\mu \| \) ≥ 0 with || \(\mu \| = 0 \) iff \(\mu = \in \),

(N3) Triangle Inquality. Yu, v & V

1 u+v1 = 1 u1 + 1 v1.

An immediate consequence of this definition together with Theorems A2,3,6 and 8 is

Theorem A 2.3, 9. An inner product space is a normed linear space under the inner product moun

 $\|\underline{u}\| := |\underline{u}| = \sqrt{\underline{u} \cdot \underline{u}}.$

INSERT material on pp. A2.3.15a, b

The inner product - magnitude inequality also can be used to define the "angle" between two elements of an inner product space. To see what to do here, we so back once again to the familian example of vectors in the Euclidean plane.

INSERT on p. A2.3.15 just after Theorem A2.3.9.

The following generalization of the triangle inequality is often useful.

Theorem A2.3, 9. Let V be a normed linear space. Then $\forall u, v \in V$

Parof, The second inequality in (i) is exactly the triangle inequality. To get the first, consider

 $\| \underline{u} \| = \| \underline{u} + \underline{v} - \underline{v} \|$ (add and subtract \underline{v})

 $= \|(\underline{u} + \underline{v}) - \underline{v}\| \quad (a \text{ sometivity})$

< | U+ V | + | - V | (triangle inequality)

= ||4+4||+||4|| (||-4|| = ||4||)

 $\| \underline{u} \| - \| \underline{v} \| \le \| \underline{u} + \underline{v} \|,$

Interchanging the roles of u and V, we obtain

11/211-11/211 = 11/2+1/21

= || u + v || (commutativity).

Thus, we have both $\| \| \| - \| \| \| \| \le \| \| \| + \| \| \|$ and $- (\| \| \| - \| \| \|) \le \| \| \| + \| \| \| \|$; where $\| \| \| \| - \| \| \| \| \| \le \| \| \| \| + \| \| \| \|$, which is the first inequality in (i).

Inequalities (ii) are astablished by replacing χ with $-\chi$ in (i) and using $\| -\chi \| = \| \chi \|$. \square

Rehum to p. A2.3.15,

In that case, the angle of between two vectors is and is was a puinitive concept; and the inner product was defined as

For U, y & Q (And, of course, the notion of engle makes no sense if one or both of the vectors is the zero vector.), we have

$$case = \frac{\cancel{u} \cdot \cancel{v}}{|\cancel{u}| |\cancel{v}|}$$

Now the s.h.s. makes sense in any inner product space; but to be the cosine of some angle, it must be in the interval [-1,1]. This is guaranteed by the inner product - magnitude in quality. Thus, we are lead to

Definition A2.3, 4. Let U and X be two nonzero elements of an inner product space. Then the angle Θ between them is defined by $\cos \Theta = \frac{U \cdot V}{|U \cdot V|}, \quad 0 \le \Theta \le \pi .$

Exercise A 2.3.10. Let a and by be nonzero elements of any inner product space. Let Θ be the angle letween them, and write C = Q - b. Establish the law of cosines:

|C|2= |a|2+|b|2-2|a||b|ασθ.

While it is amusing to be able to speak of the angle between two elements of spaces such as C°([a,b]) where we usually do not think geometrically, it is not of much utility except for suggesting the definitions of "orthogonal" $(\theta = \pi/2)$, "parallel" $(\theta = 0 \text{ or } \pi)$, The found definitions below are not given Livetly in terms of Q, because this would restrict us to nonzero elements, and "same Lirection" (0=0).

Definition A2,3,5, Two elements it and & of an inner product space are orthogonal if U.V = 0.

Exercise A2.3,11. Let a and b be orthogonal elements of any inner product space, and set S= a-b. Prove the Pythogovan Theorem: \$ S=0-5

101= 1012+1212,

Continuing in this same line of thought, parallel should mean 0=0 or TT, which implies cos 0 = ±1. Then the definition of angle gives u. x = ± | u | | x |, which corresponds to equality in the inner product - magnitude inequality. But we can do better than this if we take a slightly different line of thought, For vectors in the Euclidean plane, parallel vectors lie along the same line; i.e., are is a scalar multiple

Sovetimes, plupandicular.

ites the advantage of not being restricted to inner product spaces.

Space (not necessarily on inverpendent, u and v of any linear if the set {u, x} is linearly dependent,

In view of the preceding remarks, the following result is not surprising.

Theorem A2.3.11. Two elements u and v of an inner product space are parallel iff $u \cdot v = \pm |u| |v|$.

Proof. Exercise A2.3.12. (Hint: Parallel \Leftrightarrow linear dependence $\Leftrightarrow \exists x, \beta \text{ not both zero } \ni x u + \beta v = Q$. In the difficult direction $(u \cdot v = \pm |u||v|) \Rightarrow \text{parallel})$, use the positive definiteness, $(x u + \beta v) \cdot (x u + \mu v) \ge 0$ with = 0 iff $x u + \beta v = 0$, together with an intelligant choice of x and β .)

The next theorem is intuitively obvious, even if its proof is not.

Theorem A2.3.12. Let V be an inner product space. If $V \in V$ is athogonal to all $V \in V$ which are athogonal to $W \in V$, then W is parallel to W.

¹⁰ften, collinear.

Proof. If w = Q, then there is nothing to prove because, in view of Theorem A2.2.3, all elements of V are possible to Q. . . we assume that $w \neq Q$ for the remainder of the proof.

Let ≥,B ∈ R and consider

by hypothesis. With the particular choice

$$\alpha'=1$$
 and $\beta'=-\frac{\underline{u}\cdot\underline{w}}{\underline{w}\cdot\underline{w}}$,

we also have

$$\left(\alpha' \underline{\omega} + \beta' \underline{\omega} \right) \cdot \underline{\omega} \; = \; \left(\underline{\omega} - \frac{\underline{\omega} \cdot \underline{\omega}}{\underline{\omega} \cdot \underline{\omega}} \; \underline{\omega} \right) \cdot \underline{\omega} \; = 0 \; .$$

Hence, $\lambda' w + \beta' w$ is orthogonal to w'; i.e., it is one of the χ' 's. Now return to (*) and take $\chi = \lambda' w + \beta' w$, $\alpha = \lambda'$, and $\beta = \beta'$ to get $(\lambda' w + \beta' w) \cdot (\lambda' w + \beta' w) = 0$.

Then by the positive definiteness if the inner product, $\alpha'u + \beta'vv = 0.$

: the subset {u, w} is linearly dependent; i.e., u and w are parallel. [

This proof was shown to the author by Dr. Sang Lee,

Of cause, for two elements u and v to have the same direction should mean that the angle Θ between them be zero, Then

COSB = (4.11/1)

leads us to

Definition A2.3.7. Two elements u and v of an inner product space have the same direction if u.v = |u||x|.

Obviously, two elements that have the same direction are necessarily parallel, and the zero element has the same direction as every other element of the underlying inner product space.

The following necessary and sufficient condition for equality of two elements of an inner product = pace agrees with the customary definition of equality between ordinary geometrical vectors.

Theorem A2.3.13. Two elements of an inner product space are equal iff they have the same magnitude and the same direction; i.e.,

 $\mathcal{L} = \mathcal{L} \iff |\mathcal{L}| = |\mathcal{L}| \text{ and } \mathcal{L} \cdot \mathcal{L} = |\mathcal{L}||\mathcal{L}|.$

Proof. Suppose u=v. Then by substitution, |u|=|v|. Also by definition, |u|=|v|. $|u||u|=u\cdot u$;

and then by substitution this can be expressed as 16/18/ = 16.8.

> Conversely, suppose |u|= |x| and u.v = |u||v|. Consider $|\underline{u} - \underline{v}|^2 = (\underline{u} - \underline{v}) \cdot (\underline{u} - \underline{v})$ (definition of $|\cdot|$) = U·W - 2 U·V + V·V (properties of une product) = |\unu|^2 - 2\unu.\unu + |\unu|^2 (definition of |•1) $= |\mathcal{L}|^2 - 2|\mathcal{L}||\mathcal{L}| + |\mathcal{L}|^2 \quad (\mathcal{L} \cdot \mathcal{L} = |\mathcal{L}||\mathcal{L}|)$

> > $(|\underline{u}| = |\underline{v}|)$.

By positive definiteness, | w-x | = 0 > U-V = € ;

from which we conclude (Carry out the details as Exercise A2.3.13.) that $\psi = \psi$. \Box

An aspecially important application of the cancept of orthogonal elements of an inner product = pace is given in

Definition A2.3.8. A subset & U, Uz, ..., Uk } of an inner product space is orthonormal if \vert i, j = 1, 2, ..., k

$$u_i \cdot u_j = \delta_{ij} := \begin{cases} 1 & \text{if } i = j \\ 0 & \text{if } i \neq j \end{cases}.$$

Thus, the elements of an orthonormal subset are mutually orthogonal, and they are normalized in the sense that they are each of unit magnitude.

Theorem A2.3.14. An orthonormal subset of an inner product space is linearly independent,

Proof. We use Theorem AZ.Z.Z. Accordingly, suppose that

where { u, u, u, uk} is the orthonormal subset under consideration. Take the inner product of this equation with u;) 2 and conclude,

The symbol Si; so defined is Kronecker's delta.

² This, of course, is permitted, of Theorem A2.3.5.

with the did of the distributionity and homogeneity of the inner product and Theorem A2,3.3, that $\alpha_i = 0$, i = 1, 2, ..., k. \square

Since computations involving linear combinations of orthonormal subsets are partialarly simple, the following theorem and the alognithm used to prove it are of considerable significance.

Theorem A 2.3.15. Every finite-dimensimal inner product space has an orthonormal basis,

Broof. Let V_n denote any n-dimensional inner product space. By definition, V_n has a basis, say $\{\xi_1, \xi_2, \dots, \xi_n\}$. We use the Gram-Schmidt procedure to construct an arthonormal basis from $\{\xi_1, \xi_2, \dots, \xi_n\}$.

First set $\alpha_1 = c_1$, and lefine $\alpha_2 = c_2 + x \alpha_1$, where $\alpha \in \mathbb{R}$. The elements α_1 and α_2 will be arthogonal if

 $Q_1 \cdot Q_2 = Q_1 \cdot Q_2 + Q_1 \cdot Q_1 = 0$.

Since $a_1 = e_1$, $a_1 \neq 0$ by Theorem AZ, Z, 3, and then $a_1 \cdot a_2 > 0$ by the positive definiteness of the inner product, i. we can solve this equation for α . The element a_2 so constructed will be $\neq 0$ because $a_2 = 0 \Rightarrow \{e_1, e_2\}$ is linearly dependent which is ruled out by Theorem AZ, Z, 4.

Next, define $a_3 = e_3 + \beta a_1 + \delta a_2$.

as will be orthogonal to a, and as if

 $\underline{a}_{1} \cdot \underline{a}_{3} = \underline{a}_{1} \cdot \underline{e}_{3} + \beta \underline{a}_{1} \cdot \underline{a}_{1} = 0$

and

 $Q_2 \cdot Q_3 = Q_2 \cdot Q_3 + \chi Q_2 \cdot Q_2 = 0.$

Since $\alpha_1 \neq 0$ and $\alpha_2 \neq 0$, we can solve these equations for β and δ , Os above, the linear independence of $\{e_1, e_2, e_3\} \Rightarrow \alpha_3 \neq 0$.

Continuing in this fashion, we finally obtain a subset { a, a, ..., an } of n mutually orthogonal nanzero elements of Vn.

 $\left\{\frac{\underline{a}_{1}}{|\underline{a}_{1}|}, \frac{\underline{a}_{2}}{|\underline{a}_{2}|}, \dots, \frac{\underline{a}_{n}}{|\underline{a}_{n}|}\right\}$

is orthonounal. By Theorem A2.3.14, this subset is linearly independent; and hence by Theorem A2.2.8, it is a lasis for V_n . \square

Exercise A2.3,15. Consider the inner product space obtained by assigning Rn the standard inner product (see Theorem A2.3.1.). Show that the standard basis (see Exercise A2.2.6.) for Rn is orthonormal.

Generally, we have been senoting a typical lasis of an n-dimensional linear space by { £1, £2, ..., £n }. Because

Of course, a formal mathematical induction argument would be appropriate here. Work this out as Exercise A2.3,14.

of the peculiarities of orthonormal loses, the following variant of this notation will be helpful.

Definition A2.3, 9. Given an n-dimensional inver product $\frac{space}{space}$, $\frac{say}{say}$ $\frac{V_n}{1}$, $\frac{Se_{<17}}{se_{<17}}$, $\frac{e_{<27}}{se_{<17}}$, $\frac{e_{<17}}{se_{<17}}$, $\frac{e_{<17}}{se_{<17}}$ $\frac{denotes}{se_{<17}}$ and the corresponding components $\frac{s}{s}$ $\frac{denotes}{denoted}$ $\frac{denoted}{denoted}$ $\frac{denoted}{denoted}$

The computational simplifications afforded by the utilization of orthonormal bases are a result of

Theorem A2.3.16, Let $\{\xi_{217}, \xi_{227}, \dots, \xi_{2n7}\}\$ le au outhonormal basis for an n-dimensional inner product space V_n . For $u \in V_n$, write

Thou

U <i>i>= € <i>i>· \(\omega \) i = 1,2,..., n

and

W.X = U <17 V <17 + U <27 V <27 + ... + U <10 V <10>...

Proof, Exercise A2.3, 16. 1

¹ G. Theorem A2,3,2,

² See Definition AZ.Z.4.

That the equation $u^{(i)} = e_{(i)} \cdot u$ matches superscripts and subscripts is a notational defect characteristic of the use of an orthonormal basis.

Note that the inner product when written in terms of components relative to an orthonormal basis has the same special form as the standard inner product for R". In view of the isomorphism between Vn and R" given in the proof of Theorem A 2,2,13 and Exercise A 2,3,15, this is not surprising. If the inner product in terms of components has some other form, then the associated basis is not orthonormal.

^{&#}x27;G. Theorem AZ. 3.1.

Supplementary Reading

GEL'FAND, Lectures on Linear Algebra

GREUB, Linear Algebra

HALMOS, Finite-Dimensimal Vector Spaces

LICHNEROWICZ, Linear Algebra and Analysis

MARTIN and MIZEL, Introduction to Linear Algebra

MICHEL and HERGET, Mathematical Foundations in Engineering and Science

NAYLOR and SELL, Linear Operator Theory in Engineering and Science

NICKERSON, SPENCER, and STEENROD, Advanced Calculus

NOLL, Finite-Dunewional Spaces

NOMIZU, Fundamentals of Linear Algebra

ODEN, Applied Functional Analysis

A2.4, Linear Transformation:

Functions on linear spaces into linear spaces which have the additional property that they preserve sums and scalar multiples occur so frequently that it is worthwhile to gather came of their general properties in me place. The present being section is in no sense a complete treatment of this visit topic. Some spacial cases of linear transformations will be developed more fully in later sections.

function A2.4.1. Let Il and I be real linear spaces. A function f: U -V is a linear transformation of it preserves sams and scalar multiples; i.e., if it has the following properties:

(21) Additionly.
$$\forall u, v \in \mathcal{U}$$

 $f(u+v) = f(u) + f(v)$;

(12) Homogeneity,
$$\forall x \in \mathbb{R} \text{ and } \forall y \in \mathbb{R}$$

 $f(xy) = x f(y)$.

It is important to realize that in (L1) the addition operations on the left- and right-hand sides me those of "land V, respectively; similarly, for the scalar multiplication operation in (L2).

^{&#}x27;Often, linear operator, especially when U=V.

Actually, we have already seen some examples of linear teams form ations. The general isomorphism of Definition A2.2.5 and the particular isomorphism from Un to R' developed in the perof of Theorem A2,2,13 are obviously linear transfolurations. For a slightly more subtle example, Let when a fixed element of in inner product space V and Lefine f: V > R by

$$f(\underline{u}) = \underline{u} \cdot \underline{w}$$
.

Thou

$$f(\alpha + \gamma) = (\alpha + \gamma) \cdot \alpha = \alpha \cdot \alpha + \gamma \cdot \alpha = f(\alpha) + f(\beta)$$

and

$$f(x\underline{u}) = (x\underline{u}) \cdot \underline{w} = x(\underline{u} \cdot \underline{w}) = x f(\underline{u})$$

so this f is a linear transformation on V to IR.) Since the when product is a function on VXV which is linear in either variable when the other is held fixed, it is often called a bilinear transformation.

INSERT material on p. A2.4.2a

A key property of a linear transformation is that it mays the zero element of the domain into the sero element of the domain into the sero element of the codemain.

Theorem A2.4.1. Let of be a cincar transformation on a linear space U to a linear space V. Then

In this contest, R is viewed as a linear space (of. Exercise A2.1,2).

INSERT for pA2.4.2

Differentiation and integration provide more sopinisticated examples. E.g., the mapping $f\mapsto f'$ is a linear transformation on C'([a,b]) to C'([a,b]). In the other direction, the mapping $g\mapsto h$ where h is defined by $h(\alpha) = \int\limits_{\alpha}^{\times} g(\S)d\S$ is a linear transformation on C'([a,b]) to C'([a,b]). Here, $C'([a,b]) = \{f: f \text{ is a continuous by differentiable fuc. on } [a,b] \text{ to } R\}.$

return to p. A2.4.2

f(Q) = Q.

Proof. Exercise A2.4.1. (Hint: Make intelligent choices of the variables in citter the additivity on the hanogeneity property.)

It is not difficult to see that the projecties of abbitivity and honogeneity can be combined together into a single statement.

Theorem A2.4.2. Let f be a function on a linear space Il to a linear space V. Then f is a linear transformation if

f (xx+ Bx) = xf(x)+Bf(x) Yx,BER and Yu,xeV.

Proof. Exercise A 2.4.2. D

We need some notation to save on writing.

Definition A2.4.2. Let U and V be two real linear spaces.
Then Lin (U,V) >2 is the set of all linear transformations on U to V.

the 2 on the 1. h.s. is the zero element of V.

² Sometimes, L(U,V), Hom (U,V), H(U,V).

Definition A2.4.3. The sum of two linear transformations $f_1, f_2 \in \text{Lin}(\mathcal{U}, \mathcal{V})$ is the function

$$\left(f_1 + f_2\right) : \mathcal{U} \to \mathcal{V}$$

defined by $(f_1 + f_2)(\alpha) = f_1(\alpha) + f_2(\alpha)$, $\alpha \in \mathcal{U}$.

The scalar multiple of $f \in Lin(U,V)$ by $d \in \mathbb{R}$ is the function $(af) = U \rightarrow V$

Lefined by $(\alpha f)(u) = \alpha f(u), u \in U.$

An easy but important result is that sums and scalar multiples of linear transformations are themselves linear transformations.

Theorem A2.4.3, Let $f_1, f_2 \in Lin(V, V)$. Then $(f_1 + f_2) \in Lin(V, V)$,

Let $f \in Lin(V,V)$ and $x \in \mathbb{R}$. Then

 $(\alpha f) \in Lin (u, v)$,

Proof. Exercise A2.4,3, []

of course, two linear transformations are viewed as equal of they are the same function. Formally,)

Definition A2.4.4. Two linear transformations $f_1, f_2 \in \text{Lin}(\mathcal{U}, \mathcal{V})$ are equal, and we write $f_1 = f_2$, if $f_1(u) = f_2(u) \quad \forall \ u \in \mathcal{U}$.

You peolobly sense that we are on the way to seeing that Lin (U,V) is a linear space, We still need a zero element.

Theorem A2.4.4. Let U and V be linear spaces. Define the function $\Theta: U \rightarrow V$ by

 $\mathcal{Q}(\underline{u}) = \mathcal{Q}$, $\underline{u} \in \mathcal{U}$,)²

Then $\theta \in \text{Lin}(\mathcal{U}, \mathcal{V})$, θ is called the zero linear transformation.

Proof. Exercise AZ.4.4. []

²⁹ course, the 2 here is the zero element of V.

It follows from the definition of equality that there is but one zero linear transformation, Hence, it does make sense to speak of the zero linear transformation. We cannot invoke Theorem A2.1.2 here because we do not yet know that Lin (U,V) is a linear space. There is a linear space.

Theorem A.2.4.5, With the above refinitions for addition, scalar multiplication, equality, and zero, Lin (V, V) is a linear space.

Boog. Exercise AZ. 4.5. []

In view of this result, all of the general properties of linear spaces, such as those given in Theorems A 2,1,2 - 10, apply to Lin (U,V). E.g., in the present context, Theorem A 2.1,9 says

if $f,g \in Lin(\mathcal{U},\mathcal{V})$ and $\alpha,\beta \in \mathbb{R}$, then $\alpha = \beta$ and $f = g \Rightarrow \lambda f = \beta g$.

We shall regularly use such facts without additional comment.

Supplementary Reacting

BOWEN and WANG, Introduction to Vectors and Tensors, Vol. 1, Linear and Multilinear Algebra

GEL'FAND, Lectures on Linear Algebra

HALMOS, Finite-Dimensional Vector Spaces

MADTIN and MIZEL, Introduction to Linear Algebra

MICHEL and HERGET, Mathematical Foundations in Engineering

NAYLOR and SELL, Linear Operator Theory in Engineering and Science

NOLL, Finite-Dimensional Spaces

NOMIZU, Fundamentals of Linear Algebra

ODEN, Applied Functional Analysis

A2.5. Linear Functionals. Dual Bases. Covariant and Cartrovariant

Linear transformations on a linear space to the reals are of special importance.

Definition A2.5.1. Let V be a linear space. Then the linear space Lin (V, R) is called the algebraic Luck space of V and is denoted by V*, The elements of V* are called linear functionals.

A very useful tool for us will be

Theorem A 2.5.1. (Representation Theorem for Linear Functionals) Let $f \in V_n^* = \text{Lin}(V_n, \mathbb{R})$ where V_n is an n-deinensianal inner product space. Then \exists a unique element $f \in V_n$ \ni

fin) = f.u YueVn.

[&]quot;See Definition A2.4.7 and Theorem A2.4.5. Here, of course, we are viewing IR as a linear space (See Exercise A2.1.2.).

Often, simply, Sual space. The topological dual space is defined the same way except that the linear functionals are required to be brunded (see Definition 43.4.1); this is automatic when V is finite-dimensional

⁽cont.)

Note. By theorem A2.3.14, V_n has an orthonormal basis, say $\{e_{217}, e_{227}, \cdots, e_{2n7}\}$. Using the notation introduced in Definition A2.3.9, for any $u \in V_n$, we have

$$f(u) = f(u^{2/2}e_{217} + u^{227}e_{227} + \cdots + u^{2n7}e_{2n7})$$

$$= u^{217}f(e_{217}) + u^{227}f(e_{227}) + \cdots + u^{2n7}f(e_{2n7}) \quad (linearity)$$

$$= f \cdot u \quad (Theorem A2, 3, 15),$$

$$\xi = f(e_{217})e_{217} + f(e_{227})e_{227} + \cdots + f(e_{2n7})e_{2n7},$$

Here, we have taken advantage of the fact that no matter what the inner product is, it has the standard form in terms of components relative to an arthmormal casis.

Now suppose that f' is another element of V_n with this property. Then $f(u) = f \cdot u = f \cdot u$ \Rightarrow

 $(f'-f)\cdot w = 0.$

⁽cont.) 4
The infinite-dimensional counterpart is the famous Riesz
representation therein for a bounded linear functional on a Hilbert space.

⁵ g. Theorem A2,3,2,

Since this must hold $\forall \ u \in V_n$, Theorem $A2,3,4 \Rightarrow f' = f$, i.e., the special element of V_n in the representation is unique. \square

The Just space turns out to have the same dimension as the original space.

Theorem A 2.5.2.) Let V_n be an n-dimensional linear space. Then the Lual space $V_n^* := \text{Lin}(V_n, R)$ is also N-dimensional. In fact, if $\{\mathcal{L}_1, \mathcal{L}_2, \dots, \mathcal{L}_n\}$ is a basis for V_n , so that $u \in V_n$ has the unique expresentation

u=u'e, +u'e, + ···+ u"en,)'

Hen the subset of linear functionals $\{f', f', ..., f''\} \subset V_n^*$ Lefined by $f^i(u) = u^i$, $u \in V_n$, i = 1, 2, ..., n

is a lasis for Vn.

Proof. By Definition A2.2.3, we merely need to show that Ef', f^2, \dots, f^n } is a lasis for V_n^* . That the function $f^i: V_n \to \mathbb{R}$ defined by

 $f^{i}(u) = u^{i}$, $u \in V_n$

of Definition AZ. Z. 4 and Theorem AZ. Z. 11.

Mally is a linear functional is left as Exercise A2.5.1. Let
$$f$$
 be an arbitrary element of V_n^* . Then for any $u \in V_n$

$$f(u) = f(u'e_1 + u'e_2 + \cdots + u^n e_n) \quad (substitution)$$

$$= u'f(e_1) + u^2f(e_2) + \cdots + u^n f(e_n) \quad (linearity)$$

$$= f(e_1)f'(u) + f(e_2)f'(u) + \cdots + f(e_n)f'(u) \quad (substitution).$$
Thus, the subset $\{f', f^2, \dots, f^n\}$ spans V^* . To establish linear independence, we use Theorem A2.7.2. Accordingly, consider

consider $\alpha_1 f'(\underline{u}) + \alpha_2 f^2(\underline{u}) + \cdots + \alpha_n f''(\underline{u}) = 0.$

When U=ek, this becomes

But by the definition of fi

$$f'(e_k) = (e_k)^i = S_k^i, \quad)^2$$

so that

:. $\{f', f'', \dots, f''\}$ is linearly independent; and by Definition A2.7.2 it is a losis for V'''. \square [1] See Exercise A1.2.3 for remorts on the addition and 2.3.8 in Definition 2.3.8, $S_{k}^{i} := \{0 \text{ if } i = k \}$ functions.

We had no need of an inner product in Theorem A 2.5.2; But since "In is finite-dimensional, we can always equip it with an inner product. Consider this done. Now, the elements of the basis $\{f', f', ..., f''\}$ for the Junal space "In" are linear functionals. ... by the representation theorem \exists a unique $e' \in V_n \ni$

fi(w) = ei. w Yuev,

On remembering that f^{i} is defined by $f^{i}(\alpha) = \alpha^{i}$,

we are lead to

Theorem A2.5,3, Let V_n be an n-dimensional inner product space. Given a basis $\{E_1, E_2, \dots, E_n\}$ for V_n , so that $u \in V_n$ has the unique representation

u= u'e, + u'e, + ... + u'en,

 $\exists \underline{a} \text{ unique subset} \{\underline{e}^1, \underline{e}^2, ..., \underline{e}^n\} \subset V_n \ni \underline{e}^i \cdot \underline{u} = \underline{u}^i \quad \forall \ \underline{u} \in V_n.$

The subset {e', e', ... en} is a basis for Vn. It is called

¹ Theorem AZ. 5. 1,

the dual Jasis' of the original basis { e, e2, -, en }.

Proof, All that remains to be established is that {e', e', ..., e''} is a lasis. In view of Theorem A 2.2.8, it suffices to show that {e', e^2, ..., e''} is linearly independent. As usual, we use Theorem A 2.2.2. Accordingly, consider

a, e'+ a e2+ --- + x, en = Q.

Take the inner product of this equation with any $u \in V_n$ and use the linearity of the inner product to get $A, \stackrel{?}{\sim} u + A_2 \stackrel{?}{\sim} u + \dots + A_n \stackrel{?}{\sim} u = 0 \cdot u = 0$.

But ein = ui, so this becomes

 $(4) \quad \propto_{n} u' + \chi_{n} u^{2} + \cdots + \chi_{n} u^{n} = 0.$

Choose $w = e_k$, so that $u^i = (e_k)^i = S_k^i$. Then $(x) \Rightarrow x_1 S_k^1 + \alpha_2 S_k^2 + \cdots + \alpha_n S_k^n = \alpha_k = 0$,

and {e', e2, ..., en} is linearly independent. [

Let us review what we have done here. We started with a Basis { e, e, e, e, for an 11-dimensional inner product spore ",, and used it to construct a Basis { f, f, ..., f ? for the

Sometimes, reciprocal basis.

(algebraic) dual space V_n^* . Then the representation theorem for linear functionals was applied to each of the f^i :

$$f^{i}(\alpha) = e^{i} \cdot \alpha , \alpha \in V_{n}$$

The subset $\{e^1, e^2, \dots, e^n\} \subset V_n$ thus obtained is the dual basis of the primal basis $\{e_1, e_2, \dots, e_n\}$. Note that the dual basis is a basis for V_n — not for the dual space V_n^* .

The following theorem explains the terminology "suiprocal basis" which is cometimes used in place of Lucil Casis.

Theorem A2.5.4. Let $\{\xi_1, \xi_2, \dots, \xi_n\}$ and $\{\xi', \xi_2, \dots, \xi''\}$ be a basis and its dual for an n-dimensional inner product \underline{spaic} . Then

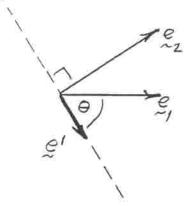
$$e^{i} \cdot e_{j} = \delta_{j}^{i}$$
, $i,j = 1,2,...,n$.

Proof. By Theorem A2.5.3, $\forall u$ in the underlying space $e^{i} \cdot u = u^{i}$.

Choose
$$u = e_j$$
 to get
$$e^{i} \cdot e_j = (e_j)^{i} = S_j^{i},$$

in which the last step follows from

We can shoot some light an $e^i \cdot e_j = e^i$ by considering the linear space of vectors associated with the <u>Euclidean plane</u>, Supprose that $i \in [0, e_2]$ is a given basis for this space, We wish to find the dual basis $i \in [0, e_2]$. Since $e^i \cdot e_j = e^i$ is orthogonal to e^i . Hurs, e^i lies along the Jotted line.



We also have $c' \cdot c_1 = S_1' = 1$. But $c' \cdot c_1 = |c'||e_1|\cos\theta$, where θ is the angle between c' and c_2 , ... $\cos\theta > 0$, so θ must be the acute angle insticuted in the figure; i.e., c' must point as shown. Now that θ is known, the length of c' is determined by $1 = |c'||e_1|\cos\theta$. The vector c' can be obtained in the same manner.

Since no ambiguities should up in the above construction, if suggests that the conditions $e^i e_j = 8^i_j$ uniquely characterize the Sual Basis. This is, indeed, the case, formally, we have

Theorem A 2.5.5. Let $\{e_1, e_2, \dots, e_n\}$ be a basis for an n-dimensional inner product space v_n . If the subset $\{e_1', e_2', \dots, e_n'\} \subset v_n$ has the projecty that

 $e^{i} \cdot e_{j} = s_{j}^{i}$, i, j = 1, 2, ..., n,

then { e', e2, ..., en} is the dual basis of { e, e, ..., en}.

Proof. Let $u \in V_n$ so that

W = U1 E, + U2 Ez + ... + Un En.

Take the inner product of this equation with ei and use the linearity of the inner product to get

 $e^{i} \cdot \mathcal{U} = u'e^{i} \cdot e_{j} + u^{2}e^{i} \cdot e_{z} + \cdots + u^{n}e^{i} \cdot e_{n}$ $= u^{i} \left(e^{i} \cdot e_{j} = \delta_{\delta}^{i} \right).$

Then by Theorem AZ. S. 3, $\{e', e^2, \dots, e^n\}$ is the dual Basis of $\{e_1, e_2, \dots, e_n\}$, \square

An immediate corollary is

Theorem A2.5,6, Let & e, e, e, e, i and & e'e, e, e, en's le a casis and its dual for an n-dimensional inner product spare. Then the dual Basis of the dual Basis & e', e', ..., en's is the original basis & e, e, e, en's.

A slightly more difficult corollary of Theorem A2.5.5 is

Theorem A2.5.7. Let {e, e2, ..., en } and {e!e2, ..., en}
be a basis and its heal for an n-dinewional inner product

space, Then {e, e2, ..., en } is orthonormal if

 $e^{i} = e_{i}$, $i = 1, 2, \dots, n$.

Proof. Exercise AZ.S.Z. [

Now we are in position for more -tandard terminology; however, it would be inaccurate to describe our approach to this topic as standard.

Definition A 2.5.2. Let $\{e_1, e_2, \cdots, e_n\}$ and $\{e', e^2, \cdots, e^n\}$ be a basis and its dual for an n-dimensional inner product e^n so that e^n is the unique representations

u=u'e, +u'e, +-- + une,

and

 $\mathcal{L} = \mathcal{L}_{1} e^{1} + \mathcal{L}_{2} e^{2} + \cdots + \mathcal{L}_{n} e^{n}.$

The scalars u^i (the components if u w.r.t. to the original basis $\{e_i, e_2, \dots, e_n\}$) are the contravariant components if u, while the scalars u_i (the components of u w.r.t. the dual basis $\{e', e^2, \dots, e^n\}$) are the covariant components of u.

See Definition AZ. 3.10.

The mnemonic "co goes below" is aseful in helping this ferminology straight. There is nothing special about the original basis often than that it is the one with which we happen to start.

Theorem A2.5.8. Let $\{\xi_1, \xi_1, \dots, \xi_n\}$ and $\{\xi_1, \xi_2, \dots, \xi_n\}$ be a losis and its dual for an n-dimensional inner product space \mathcal{V}_n .

Then for $u \in \mathcal{V}_n$ $u^i = \xi^i \cdot u \quad , \quad i = 1, 2, \dots, 11,$ and $u_i^i = \xi_i \cdot u \quad , \quad i = 1, 2, \dots, n.$

Proof. The statement $u^i = e^i \cdot u$ is the defining property of e^i . To establish the second statement, consider $u = u_i e^i + u_i e^2 + \cdots + u_i e^{ii}$.

Take the inner product of this squation with \underline{e}_i to get $\underline{e}_i \cdot \underline{u} = \underline{u}, \underline{e}_i \cdot \underline{e}' + \underline{u} \underline{e}_i \cdot \underline{e}^2 + \dots + \underline{u}_n \underline{e}_i \cdot \underline{e}^n$ (linearly of the wind product) $= \underline{u}_i \quad (\underline{e}_i \cdot \underline{e})' = \underline{S}_i \quad . \quad \square$

As an immediate consequence of the last two theorems, we have

Theorem A2.5,9. Let Un be an ordinavioural inner product space, Then relative to an orthonormal basis for Un, the

covariant and contravariant components of any element $u \in V_n$ coalesce; i.e., $u_{\langle i \rangle} = u^{\langle i \rangle}$, $i = 1, 2, \cdots$, n.

That the equation above matches subscripts and superscripts is a notational legent that will manifest itself whenever an orthonormal lasis is employed.

It will be convenient to adopt the

Summation Convention. If an unspecified index appears exactly twice in the same monomial, once as a subscript and once as a superscript, then summation from 1 to no over this index is implied. It is the dimension of the underlying linear = pace.

For example,

 $u_i e^{\bar{\iota}} = \sum_{i=1}^n u_i e^{\bar{\iota}} = a_i e' + u_i e' + \cdots + u_n e''.$

The following notation will also be useful.

Definition A2.5,3, Let {e,e,, en} and {e',e,, e"} be a basis and its Lual for an 1-dimensional inner product = pare Un. Then the associated g-symbols are

¹ i, j, etc. as opposed to 1,2, etc.

$$q_{ij} = e_i \cdot e_j$$
, $q^{ij} = e^i \cdot e^j$, $q_j^i = e^i \cdot e_j$, $i, j = 1, 2, \dots, n$.

Some of the more important properties of the g-symbols are given by

Theorem AZ. 5, 10, Refer to the notation of Definitions A2.5,2 and 3

Then $(i) \ g_{ii} = g_{ij}, \ g^{ii} = g^{ij}, \ g_{i} = g_{i}^{i}, \ i,j = 1,2,...,n_{j}$

(ii) gi = 8i, i, j = 1,2, -, n;

(iii) For any u, v & Vn,

 $u \cdot v = g_{ij} u^i v \dot{s} = g^i \dot{u}_i \dot{v}_j = g^i \dot{u}_i v \dot{s} = g^i \dot{u} \dot{s} \dot{v}_i = u_i \dot{v}_i = u^i \dot{v}_i$

(iv) For any u & Vn,

 $|u| = \sqrt{g_{ij}u^{i}u\dot{s}} = \sqrt{g^{i}\dot{s}u_{i}u_{j}} = \sqrt{u_{i}u^{i}}$

(v) qingki = Si , gikgn = Si , i,j=1,2,...,n;)1

(vi) The motrices [qi;] and [gis] are positive definite;

(Vii) ded[qis] = (ded[qis])".

Thus, the matrix [gis] is the inverse of the matrix [gis]. Here, [gis] stands for the matrix whose ij-element is gis.

Plot. (i):
$$g_{ji} = e_{j} e_{i}$$
 (lefinition)

$$= e_{i} \cdot e_{j}$$
 (commutationity of the inner percent)

$$= g_{ij} \cdot (definition).$$

Similarly, $g_{ij} = g_{ij} \cdot (definition)$

$$= g_{ij} \cdot (definition)$$

$$= g_{ij} \cdot (definition)$$

$$= g_{ij} \cdot (definition)$$

$$= e_{i} \cdot e_{i} \cdot (definition).$$

(ii): Contained in the people of (i).

(iii): $u \cdot v = (u_{i} e_{i}) \cdot (v_{i} e_{j}) \cdot (component sepassarktion)$

$$= u_{i} v_{i} \cdot (e_{i} \cdot e_{j}) \cdot (linearity of inner peoduct)$$

$$= g_{ij} u_{i} v_{i} \cdot (definition).$$

Similarly, $u \cdot v = g_{i} \cdot u_{i} \cdot v_{i} \cdot v_{i}$, and

$$\mathcal{U} \cdot \mathcal{V} = g_i^{\delta} u^i v_j$$

$$= \mathcal{S}_i^{\delta} u^i v_j \quad ((ii))$$

$$= u^{\delta} v_j \quad (\text{property of Kronecker': Jella}).$$

$$Similarly, \quad u \cdot v = u_i v^i,$$

$$(iv): \quad |u|^2 = u \cdot u \quad (\text{Jefinition})$$

$$= g_{ij} u^i u_j = g^{ij} u_i u_j = u_i u^i \quad ((iii)).$$

$$(v): \quad g_{ik} g^k s = (e_i \cdot e_k) (e^k \cdot e^k) \quad (\text{Jefinition})$$

$$= (e_i)_k (e_i^{\delta})^k \quad (\text{Theorem A2.5.8})$$

$$= e_i \cdot e_i^{\delta} \quad ((iii))$$

$$= \mathcal{S}_i^{\delta} \quad (\text{Theorem A2.5.4}).$$

Similarly, gikgnj = si.

(Vi): Consider $u \cdot u = g_{ij} u^i u^j$, By the positive definiteness of the inner product $u \cdot u \geq 0$ with = 0 if u = Q. By Theorem A2.2.12, u = Q iff $u^i = 0$, $i = 1, 2, \cdots, n$. Hence, the quadratic form $g_{ij} \cdot u^i u \bar{s}$ and the associated matrix $[g_{ij}]$ are positive definite by definition, Similarly, for $[g^i \bar{s}]$.

(vii):
$$S_i^{\dot{g}} = g_{ik} g^{k\dot{g}}$$
 ((v))
 $\Rightarrow \det [S_i^{\dot{g}}] = \det [g_{ik} g^{k\dot{g}}]$
 $= \det [g_{ik}] \det [g^{lm}]$ (For equare matrices A and B, $\det(AB) = \det(A) \det(AB)$).

so $\det[S_i^i] = 1$ Let $[g_{ik}] \det[g^{\ell m}] = 1$. Since the product on the l.l.s. equals 1, neither factor can vanish; and we can write $\det[g_{ik}] = (\det[g^{\ell m}])^{-1}. \square$

Our next result shows how the g-symbols ran be used to convert between covariant and contravariant components.

Theorem A2.5, 11, Refer to the notation of Definitions A2.5. 2 and 3. Then for any $u \in V_n$,

$$u^{i} = g^{ij}u_{j}$$
 and $u_{i} = g_{ij}u\dot{\delta}$, $i = 1, 2, \dots, n$.

Prof.
$$u^{\xi} = e^{\xi} \cdot u$$
 (Theorem A2.5.8)

$$= e^{i} \cdot (u, e^{j}) \qquad (component inpresentation)$$

$$= u_{j} \cdot (e^{i} \cdot e^{j}) \qquad (linearity of the inner product)$$

$$= g^{ij}u_{j} \qquad (lefinition of g^{ij}),$$

Similarly, U= qi, ui. 1

For obvious reasons the operations indicated in Theorem A2.5,11 are often referred to as "raising and lowering indices". The next result shows that we can do the same thing with the basis elements.

Theorem A2.5.12, Refer to the notation of Definition A2.5.3, Then

 $e^{i} = g^{ij}e_{j}$ and $e_{i} = g_{ij}e^{j}$, i=1,2,...,n,

Proof. Exercise A7.5.3. \square

Supplementary Reading

Same as for \$ A2,3 plus

COLEMAN, MARKOVITZ, and NOLL, Viscometric Flows of Non-Newtonian Fluids

(Appendix on Mathematical Concepts)

LICHNEROWICZ, Tensor Calculus

A2.6. Transformations of Bases

In this section, we investigate how the components of an element of a finite-dimensional linear space change when the basis is changed. We start by considering the relationships that must hold between my two loses.

Theorem A2.6.1. Let { e, e, e, e, e, } { e', e', e, en} and { e, e, e, en} and { e, e, e, en} and their duals for an n-dimensional inner product = paic V,, Define

Then $\begin{array}{lll}
\overline{Ren} & \alpha_i^j = \overline{e}^{j \cdot e_i}, & \beta_i^j = \overline{e}^{j \cdot e_i} & (i,j=1,2,...,n). \\
\overline{e}_i = \alpha_i^j \overline{e}_j, & \underline{e}^i = \beta_i^j \overline{e}^{j} & (i=1,2,...,n), \\
\overline{e}_i = \beta_i^j e_j, & \overline{e}^i = \alpha_i^j \underline{e}^{j} & (i=1,2,...,n), \\
\underline{and} & \alpha_i^j \beta_k^k = \delta_i^k, & \alpha_i^j \beta_k^l = \delta_k^i & (i,k=1,2,...,n).
\end{array}$

Ploof. By Definitions A2,2,2 and 3, ϵ_i can be written as a linear combination of the wembers of the Basis $\{\bar{e}_i, \bar{e}_i, \bar{e}_i, \bar{e}_n\}$; i.e.,

e = x ! E;.

See Theorem A2,5,3.

² g. Theorem A2.3.2.

By Theorem AZ. 5. 8, the components of are given by dd = e. · ē d ei= sied with Bi=ei.ej. With the above definitions for the x's and the B's, $\alpha \stackrel{\ell}{:} \beta \stackrel{k}{!} = (e \stackrel{\bar{e}}{:} e \stackrel{\bar{e}}{:} e) (e^{k} \cdot \bar{e}_{\ell})$ = (e;) (ek) = (Theorem A2.5,8) = e. ek (Theorem A2.5.10 (iii)) = 8 k (Theorem A2.5.4). Similarly, x1 3 = 5 ,

Finally, we can use the above results to write $\beta_{i}^{j} \in_{j} = \beta_{i}^{j} (A_{j}^{k} = A_{k}) = S_{i}^{k} = A_{k} = A_{i}^{k}$

Here, (Ei) I denotes the lth continuous temperant of Ei w.r.t. the Barred Bosis. Marking the index (nother than the perrel letter as in Treorem A2.6,2) to indicate the underlying Bosis is occasionally advantageous (cf. LICHNEROWICZ's Tensor (alculus).

$$\alpha_{j}^{i} = \overline{e}_{i}$$
. \square

As an important consequence of this Henrem, we have

Theorem A2.6.2, Let the hypotheses of Theorem A2.6.1 hold. Let $u \in V_n$ so that

$$u = u^i e_i = u_i e^i = \overline{u}^i \overline{e}_i = \overline{u}_i \overline{e}^i$$

Then, for i=1,2, ..., 17,

u'= x' u' -- transformation rule for contravariant components,

Ui = Bi U --- transformation rule for covariant components,

and

$$u^{i} = \beta^{i} \overline{u}^{j}$$

Proof. We use the soults of Theorem A2.6.1 tegether with the uniqueness of components to write

$$\mu = \mu^i e_i = \mu^i \chi^j \epsilon_j = \overline{\mu}^j \epsilon_j \Rightarrow \overline{\mu}^j = \chi^j \mu^i,$$

$$\mathcal{L} = \mathcal{U}_i \mathcal{L}^i = \mathcal{U}_i \mathcal{L}^i \bar{e}^j = \bar{\mathcal{U}}_i \bar{e}^j = \bar{\mathcal{U}}_i \bar{e}^j \Rightarrow \bar{\mathcal{U}}_i = \mathcal{L}^i \mathcal{L}^i .$$

¹ Theorem A2.2.11,

Then
$$\beta_{j}^{k} \bar{u}^{j} = \beta_{j}^{k} (\alpha_{j}^{j} u^{i}) = \delta_{i}^{k} u^{i} = u^{k},$$

$$\alpha_{k}^{j} \bar{u}^{j} = \beta_{j}^{k} (\beta_{j}^{j} u^{i}) = \delta_{k}^{i} u^{i} = u_{k}.$$

Obviously, if we have a list of reals of length n, say $(u', u^2, ---, u^n)$, we can generate an element, say u, of v_n by letting the u' be the components of u w.r.t. some basis, say $\{e_1, e_2, ---, e_n\}$:

$$\omega := \omega^i e_i$$
.

A common occurrence in the physical sciences is that by some natural scheme we have real lists associated with each Gasis, say

$$(u',u',--,u'')$$
 with $\{\underline{e}_{n},\underline{e}_{2},--,\underline{e}_{n}\}$, $(\bar{u}',\bar{u}^{2},--,\bar{u}'')$ with $\{\underline{\bar{e}}_{n},\underline{\bar{e}}_{2},--,\underline{\bar{e}}_{n}\}$, etc.

This allows us to generate an infinity of elements of Vn: $\mathcal{U} = \mathcal{U}_{i}^{e}$, $\overline{\mathcal{U}} = \overline{\mathcal{U}}_{i}^{e}$, etc.

By Theorem A2.6.3, if these elements of Vn are all the same

(i.e., if $u = \overline{u}$, etc.), then the u^i , \overline{u}^i , etc. satisfy the

 $\bar{u}^i = \lambda^i u \delta$, $u^i = \beta^i \bar{u}^j$, etc.

This leads us to

Theorem A2.6.3. Let the hypotheses of Theorem A2.6.1 hold. Let there be real list of length n associated with each Basis for v_n : $(u', u^2, ..., u^n) \text{ with } \{\underline{e}_1, \underline{e}_2, ..., \underline{e}_n\},$ $(\bar{u}', \bar{u}^2, ..., \bar{u}^n) \text{ with } \{\underline{e}_1, \underline{e}_2, ..., \underline{e}_n\},$ $(\bar{u}', \bar{u}^2, ..., \bar{u}^n) \text{ with } \{\underline{e}_1, \underline{e}_2, ..., \underline{e}_n\},$ $\underline{etc},$

The entries in these lists are the contravariant components of an elemen $u = u^i e_i = \bar{\alpha}^i \bar{e}_i = --$

of V_n iff they satisfy the transformation rules $\bar{u}^i = \alpha^i u^j$, $u^i = \beta^i \bar{u}^j$,

Obviously, there is a strictly analogous result for covariant component.

Part. Only the "if" portion remains to be proven, Accordingly, define

 $\mathcal{U} = \alpha' \underline{e}_i$, $\overline{\mathcal{U}} = \overline{\alpha}' \underline{e}_i$, etc.,

where the components ui, ui, etc. west the stated transformation rules. Then

 $\overline{U} = \overline{u}^{i} \overline{e}_{i} \quad (construction)$ $= (d_{j}^{i} u \dot{d}) \overline{e}_{i} \quad (thansformation rule)$ $= u \dot{d} (d_{j}^{i} \overline{e}_{i}) \quad (linear space axioms)$ $= u \dot{e}_{j} \quad (Theorem A 2.6.1)$ $= u \quad (construction).$

Thus, if the transformation rules me satisfied, the natural constructions all generate the same element of Vn. []

Usually in the physical sciences, elements of finite-dimensional linear spaces are defined through their components together with the requirement that their components transform according to the rules of theorem A2.6.2 under a change of basis. We shall refer to such a scheme as the component-transformation rule approach.

The next therem shows that the g-symbols obey Similar transformation rules,

Theorem A2.6.A. Let the hypotheses of Theorem A2.6.1 hold. Define

 $\frac{g'j = e^{i} \cdot e^{j}}{g'j = e^{i} \cdot e^{j}}, \quad \overline{g'j} = \overline{e^{i} \cdot \overline{e}^{j}}, \quad \underbrace{e^{j}e^{i}}, \quad (i,j = 1,2,...,n).$ $\frac{f^{n}}{g'j} = \alpha_{k}^{i} \alpha_{k}^{j} g^{kl}, \quad \overline{g}_{ij} = \beta_{i}^{k} \beta_{j}^{l} g_{kl}$

 $\frac{g}{g} \dot{s} = \lambda_{k}^{i} \beta_{j}^{i} g_{k}^{k} = S_{j}^{i} = g_{j}^{i}$ $\frac{and}{g} \dot{s} = \beta_{k}^{i} \beta_{k}^{j} \bar{g}^{kl}, \quad g_{ij} = \lambda_{i}^{k} \lambda_{j}^{l} \bar{g}^{kl},$

1 See Definition A2,5,3,

Proof. Exercise AZ. G. I. [

Note that if both of the Bases $\{\underline{e}_1,\underline{e}_2,...,\underline{e}_n\}$ and $\{\underline{e}_1,\underline{e}_2,...,\underline{e}_n\}$ are orthonormal, then the above result yields orthogonality conditions such as

for the transfounction coefficients, Such conditions are protobly familian to you from "Cartesian tensor alsobra" (usually formulated via the component-transformation rule approach) where only orthonormal bases are comployed.

Exercise A2.6.2. In the component-transformation whe approach, the inner product of two elements would be defined in terms of components as

If course, in this approach me needs to be concerned about the definition being tied fundamentally to the Basis used in the liquidion. Use the transformation rules to show that

If course, in our Gasis-fue approach to the inner product such questions never arise; and

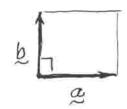
follows directly from Theorem A 2.5.10 (iii).

Supplementary Reading Same as for \$ A2,5

A 2.7. Volume Orientation Functionals. Oriented Spaces

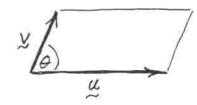
Here, we put in place the foundation for consideration of the "vector product" of two elements of a three-dimensional space in § A2.8 and the "leterminant" of a "tensor" on an n-chimensional space in § A3,5. For motivation, we revert again to vectors in the Euclidean plane,

If a and & are two orthogonal vectors in the plane, they generate a rectangle.

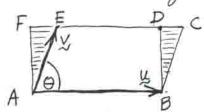


By definition, the area of a redangle is the product of the base and the reight; i.e., area {a,b} := |a||b|.

In the general case, two vectors & and & in the plane generate a parallelogram.



The area of a parallelogram, is still the product of the base and the height as is seen from the following figure.



The Base times the height gives the area of the rectangle ABDF. This area includes the area of the triangle AFE, which is not part of the area of the parallelogram generated by y and y; but it leaves out the area of the triangle BDC, which is included in the area of the parallelogram. However, the areas of these two triangles me equal. Thus, since the height is 14,15in0,

ava { \(\omega, \omega \) = \(|\omega | |\omega | |\omega | \omega \) sin \(\omega \).

Sine, by convention, the angle between two vectors is taken to be the smalle of the two possibilities, this general formula provinces a nonnegative area. It also gives the covert result in the rectangular case,

The area of a parallelegram has several injuntant properties, First, we note that if x>0 so that xy and y point in the same sinction, then

area $\{ \propto u, v \} = |\alpha u||v| \sin \theta = |\alpha||u||v| \sin \theta$.

If $\alpha < 0$ so that αu and u point in opposite directions, then area $\{\alpha u, \chi\} = |\alpha u| |\chi| \sin(\pi - \theta) = |\chi| |u| \sin \theta$.

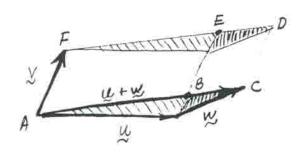
Thus, the area of a parallelogram is positive homogeneous in the sense that

area {du, y} = |d| area {u, y}.

. Similarly,

area {u, by } = 1 Blanea {u, x}.

Next, consider the area generated U+W and V, where W is another vector in the Euclidean plane. Let us exemine the case where the angle between U and W is acute.



The area generated by U+W and V is the sum of the areas of the parallelograms ABEF and BCDE, But these areas are seen to be respectively equal to the areas generated by U and V and by W and V. The case where the angle between U and W is obtuse can be randled in the same way, and thus the area of a parallelogram is additive in the sense that

area $\{ \underline{w} + \underline{w}, \underline{v} \} = \text{area} \{ \underline{w}, \underline{v} \} + \text{area} \{ \underline{w}, \underline{v} \}$.

Similarly, area $\{\underline{\omega},\underline{v}+\underline{w}\}=$ area $\{\underline{\omega},\underline{v}\}+$ area $\{\underline{\omega},\underline{w}\}$.

Finally, we note that the area of the parallelogram generated by u and x is independent of the ordering; i.e.,

ava (u, v3 = ava (v, u).

Before generalizing the notion of area to higher dimensions, we need to review the concept of a permitation.

Definition A2.7.1. If each of the first n integers $II_{n} = \{1, 2, \dots, n\}$

appears once and only once in the list (J, Jz, Jn), then (Ji, Jz, Jn) is a permutation of In. When in a permutation an integer precedes a smaller integer, the permutation is said to contain an inversion. A permutation is cure or odd according as the total number of inversions it contains is even or odd.

of course, the total number of inversions in a plumutation can be found by counting the number of smaller integers following each integer of the permutation, E.g., the permutation (6,1,4,3,2,5) contains eight inversions since

4 is followed by 1,4, 3, 2, and 5, 4 is followed by 3 and 2, 3 is followed by 2.

See HOHN's <u>Elementary Matrix Algebra</u> for a comprehensive development along the line adopted here. See NOLL's <u>Finite-Dimensional</u> Spaces for a more abstract treatment.

Definition A2,7,2. The n-dimensional alternating symbol $\varepsilon_{\sigma_1\sigma_2}$ on is defined by

$$\mathcal{E}_{\overline{J_1}\overline{J_2}\cdots\overline{J_n}} = \begin{cases} 1 & \text{if } (\overline{J_1},\overline{J_2},\cdots,\overline{J_n}) \text{ is an even permutation of } \overline{I_n} \\ -1 & \text{" "odd " " " " } \\ 0 & \text{" "not a " " " " } \end{cases}$$

The next step in our generalization of parallelogram area is

Definition A2.7.3. A volume-orientation functional Δ for an n-dimensional inner product = pace! V_n is a function on $V_n \times V_n \times \cdots \times V_n$ to R which obeys the following axioms:

(01) Multilinearity. ∆ is linear in each argument; i.e., ∀ u, u, ..., un, x ∈ Vn and ∀ α ∈ R

 $\Delta(\underline{u}_{i}, \dots, \underline{u}_{i-1}, \alpha \underline{u}_{i}, \underline{u}_{i-1}, \dots, \underline{u}_{n})$

 $=\alpha\Delta(\underline{u}_{1},\cdots,\underline{u}_{i-1},\underline{u}_{i},\underline{u}_{i+1},\cdots,\underline{u}_{n})$

and

 $\Delta(\underline{u}_{i}, \dots, \underline{u}_{i-1}, \underline{u}_{i} + \underline{v}, \underline{u}_{i+1}, \dots, \underline{u}_{n})$

 $=\Delta(\underline{u}_1,\cdots,\underline{u}_{i-1},\underline{u}_i,\underline{u}_{i+1},\cdots,\underline{u}_n)+\Delta(\underline{u}_1,\cdots,\underline{u}_{i-1},\underline{v},\underline{u}_{i+1},\cdots,\underline{u}_n)$

for any i ∈ IIn;

¹ G. Theorem A2.3.2.

(02) Skew-Symmetry. Δ is skew-symmetric w.r.t. all arguments; i.e., \forall $\sigma_1, \sigma_2, \dots, \sigma_n \in \mathbb{H}_n$

 $\Delta\left(\overset{\cdot}{\mathcal{U}}_{\sigma_{1}},\overset{\cdot}{\mathcal{U}}_{\sigma_{2}},\cdots,\overset{\cdot}{\mathcal{U}}_{\sigma_{n}}\right) = \varepsilon_{\sigma_{1}\sigma_{2}\cdots\sigma_{n}}\Delta\left(\overset{\cdot}{\mathcal{U}}_{1},\overset{\cdot}{\mathcal{U}}_{2},\cdots,\overset{\cdot}{\mathcal{U}}_{n}\right)^{1}$

(03) If Eexis, exzs, ..., exp is an orthonormal basis for Vn, then

 $\left|\Delta(\underline{e}_{\langle 1\rangle},\underline{e}_{\langle 2\rangle},...,\underline{e}_{\langle n\rangle})\right|=1.$

Next, we turn our attention to demonstrating that if is always possible to find a volume-orientation functional for Vn; in fact, there will turn out to be exactly two, with one being the negative of the other,

Theorem A2.7.1. If I a functional D: Vn x Vn x vn > R satisfying (01) and (02), it must have the form

 $\Delta(\mathcal{U}_1,\mathcal{U}_2,\cdots,\mathcal{U}_n)=\mathcal{E}_{\sigma_1\sigma_2\cdots\sigma_n}\mathcal{U}_1^{\sigma_1}\mathcal{U}_2^{\sigma_2}\cdots\mathcal{U}_n^{\sigma_n}\Delta(\mathcal{E}_1,\mathcal{E}_2,\cdots,\mathcal{E}_n),)^2$

Where { en, ez, -. , en} is any basis for Vn and

 $\mathcal{U}_i = \mathcal{U}_i \stackrel{\text{de}}{\approx}_j , i = 1, z, \dots, n.$

Here $(\sigma_1, \sigma_2, \dots, \sigma_n)$ need not be a permutation of I_n ; e.g., $(\sigma_1, \sigma_2, \sigma_3) = (3, 1, 1)$ is allowed.

2,3 As usual, the summation convention (p. A2,5,12) applies,

Then (Oz) =>

$$\Delta\left(\mathcal{U}_{1}, \mathcal{U}_{z}, \cdots, \mathcal{U}_{n} \right) = \mathcal{U}_{1}^{\sigma_{1}} \mathcal{U}_{2}^{\sigma_{z}} \cdots \mathcal{U}_{n}^{\sigma_{n}} \, \varepsilon_{\sigma_{1} \sigma_{z} \cdots \tau_{n}} \, \Delta(\varepsilon_{1}, \varepsilon_{z}, \cdots, \varepsilon_{n}). l$$

We could have worked just as well with covariant conjunents of the ci's in the above representation theorem. Then we would have obtained

$$\Delta(\underline{u}_{1},\underline{u}_{2},--,\underline{u}_{n})=\varepsilon^{T_{1}T_{2}\cdots T_{n}}u_{\tau_{1}}^{(n)}u_{\tau_{2}}^{(2)}\cdots u_{\tau_{n}}^{(n)}\Delta(\varepsilon_{1}^{\varepsilon},\varepsilon_{2}^{\varepsilon},--,\varepsilon_{1}^{\varepsilon}),$$

where ET, Jz. . . To has the same meaning as Eo, Jz ... on and

$$\mathcal{U}_i = \mathcal{U}_j^{(i)} \in \mathcal{F}_j \quad i = 1, 2, \cdots, n,$$

with {e', e', ..., e"} being the dual basis of {£1, £2, ..., £n}.

The above representations to not really lefine & because Devoluated on a basis applays on the 1. h.s.'s. We can get around this by employing an orthonormal cosis and whilizing Axiom (03).

Theolon A2.7.2. I exactly two volume-orientation functionals for a given n-dimensional inner product space Vn, and one is

the negative of the other. In fact, if { Exp, Exp, ..., Exp, 3 is an arthonounal Basis for Vn so that $u_i \in V_n$ has the unique representation

 $\frac{\mathcal{U}_{i} = u_{i}^{\langle s \rangle} e_{\langle j \rangle}, \quad i = 1, 2, \dots, n}{\Delta(u_{1}, u_{2}, \dots, u_{n}) = \pm \epsilon_{\sigma_{1} \sigma_{2} \dots \sigma_{n}} u_{1}^{\langle \sigma_{i} \rangle} u_{2}^{\langle \sigma_{i} \rangle} \dots u_{n}^{\langle \sigma_{n} \rangle},}$

with + or - according as $\Delta(e_{e1}, e_{e2}, -; e_{en}) = +1 -1.$

Part. If △ exists, Theorem A2.7.1 >>

 $\Delta(\mathcal{U}_1,\mathcal{U}_2,\cdots,\mathcal{U}_n) = \mathcal{E}_{\sigma_1\sigma_2\cdots\sigma_n} \mathcal{U}_1^{\langle \sigma_1\rangle} \mathcal{U}_2^{\langle \sigma_2\rangle} \cdots \mathcal{U}_n^{\langle \sigma_n\rangle} \Delta(\mathcal{E}_{\langle 17\rangle},\mathcal{E}_{\langle 27\rangle},\cdots,\mathcal{E}_{\langle n\rangle})$

But $(03) \Rightarrow \Delta(\underline{e}_{<1>}, \underline{e}_{<2>}, --, \underline{e}_{<3>}) = \pm 1$. Hence, if Δ exists, it must be given by either

 $\Delta(\mathcal{U}_{1},\mathcal{U}_{2},\cdots,\mathcal{U}_{n}) = + \mathcal{E}_{\sigma_{1}\sigma_{2}} \cdots \sigma_{n} \mathcal{U}_{1}^{\langle \sigma_{1} \rangle} \mathcal{U}_{2}^{\langle \sigma_{2} \rangle} \cdots \mathcal{U}_{n}^{\langle \sigma_{n} \rangle}$ $\Delta(\mathcal{U}_{1},\mathcal{U}_{2},\cdots,\mathcal{U}_{n}) = -\mathcal{E}_{\sigma_{1}\sigma_{2}} \cdots \sigma_{n} \mathcal{U}_{1}^{\langle \sigma_{1} \rangle} \mathcal{U}_{2}^{\langle \sigma_{2} \rangle} \cdots \mathcal{U}_{n}^{\langle \sigma_{n} \rangle},$

with the sign in accord with the value of Δ on the orthonormal basis employed. It is easy, and essential, to check that each of these tentative Δ 's Loes, indeed, satisfy the axioms. The details are left as Exercise A2.7.1. [The identity $\sum_{\sigma} \mathcal{E}_{\sigma, \sigma_{\sigma}, \sigma_{\sigma}} \alpha_{\tau, \sigma_{\sigma}} \alpha_{\tau, \sigma_{\sigma}} \alpha_{\tau, \sigma_{\sigma}} \alpha_{\tau, \sigma_{\sigma}} = \mathcal{E}_{\tau, \tau_{\sigma}, \sigma_{\sigma}} \det [a_{ij}]$ (see HOHN's

Elementary Matrix Algebra) will be useful in this regard.]

Note that even though the volume-orientation functional A2.7.9 was introduced abstractly through axioms expressing its A2.7.9 properties, we now know that it does exist; and we even have a formula for calculating it.

Definition A2.7.4. An n-dimensional inner product space V_n equipped with a volume-orientation functional Δ is said to be oriented by Δ .

According to Theorem A2.7.2, there are two, and only two, ways to orient V_n . For the remainder of this section V_n is an oriented n-dimensional inner product space.

Now we are in position for the following generalization of the notion of ave.

Definition A2.7.5, The volume spanned by a set of V_n elements V_n is V_n is V_n is V_n is V_n is V_n .

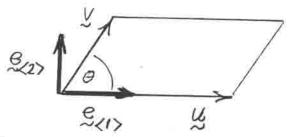
An easy but important result is

Theorem A2.7.3, The volume spanned by {u, u, u, u, cv, is independent of the choice of volume-anientation functional for vn, and also it is independent of the ordering of the n elements u, u, u, u, un.

Proof. Exercise A2.7.2. [

Exercise A2.7.3. Since V_2 models the set of vectors in the Euclidean plane, the volume spanned by a pair of monparallel elements is and x of V_2 should be the area (in the usual geometrical sense visuessed at the leginning of the present section) of the paullelogram that they

generate. Show that this is, indeed, the case. (Hint: Employ an orthonormal Basis as indicated in the figure below.)



Exercise A2,7,4, Consider the volume spanned by the singleton { U } < Vy.

I course, in Exercise A2.7.3 if U and y are parallel, then the parallelogram collepses and has zero area. In view of Definition A2.3.6 and the remarks that preceded it, this suggests

Theorem A2.7, 4, A set $\{u_1, u_2, \dots, u_n\} \subset V_n$ is linearly dependent iff $\Delta(u_1, u_2, \dots, u_n) = 0$; i.e., iff the volume spanned by $\{u_1, u_2, \dots, u_n\}$ is zero.

Epondent. Then me of the u's is a linear combination of the rest. By relabeling, if necessary, we can assume that

$$\mathcal{U}_{i} = \sum_{i=2}^{n} \alpha_{i} \mathcal{U}_{i}$$
.

Then

 $\Delta(\mathcal{U}_1,\mathcal{U}_2,\dots,\mathcal{U}_n) = \Delta(\sum_{i=2}^n \alpha_i \mathcal{U}_i,\mathcal{U}_2,\dots,\mathcal{U}_n)$ (substitution,

$$= \sum_{i=2}^{n} \alpha_i \Delta(\mathcal{U}_i, \mathcal{U}_2, \dots, \mathcal{U}_n) \quad (\text{linearity})$$

$$= \mathcal{O} \quad (\text{skew-symmetry}).$$

In the last step, we are using the fact that because of the show-symmetry, Δ vanishes whenever any two gits arguments are equal. E.g.,

 $\Delta(\underline{u}_{2},\underline{u}_{2},\underline{u}_{3},...,\underline{u}_{n}) = \varepsilon_{223\cdots n} \Delta(\underline{u}_{n},\underline{u}_{2},...,\underline{u}_{n}) ;$

but $\mathcal{E}_{223\cdots n} = 0$ since $(2, 2, 3, \cdots, n)$ is not a permutation of \mathbb{I}_n .

Conversely; suppose that $\Delta(u_1, u_2, \dots, u_n) = 0$, Assume that $\{u_1, u_2, \dots, u_n\}$ is linearly independent. Then by Theorem A2.2.8, $\{u_1, u_2, \dots, u_n\}$ is a basis for V_n . It then follows from the representation Theorem A2.7.1 that $\Delta(v_1, v_2, \dots, v_n) = 0$ for any set of n elements of v_n . This contradicts the axiom that $\Delta(e_{<1}, e_{<2}, \dots, e_{<n}) = 1$. $\{u_1, u_2, \dots, u_n\}$ cannot be linearly independent; hence, it is linearly dependent. \square

By Theorem A2.7.4, if $\{\underline{c}_1,\underline{c}_2,\dots,\underline{c}_n\}$ is a basis for V_n , then $\Delta(\underline{c}_i,\underline{c}_2,\dots,\underline{c}_n) \neq 0$. But $\Delta(\underline{c}_i,\underline{c}_2,\dots,\underline{c}_n)$ could be either positive or negative, This leads us to

Definition AZ.7.6. A basis $\{e_1, e_2, \dots, e_n\}$ for V_n is positive or negative according as $\Delta(e_1, e_2, \dots, e_n) > 0$ or

<0. Two bases { e, e, ..., en } and { \vec{e}_1, \vec{e}_2, ..., \vec{e}_n} are

 $\Delta(\underline{e}_1,\underline{e}_2,\dots,\underline{e}_n)\Delta(\underline{\overline{e}}_1,\underline{\overline{e}}_2,\dots,\underline{\overline{e}}_n)>0.$

As an immediate consequence of theorem AZ. T. Z, we have

Theorem A2.7,5. The property of like-handedness is independent of the choice of volume-orientation functional.

We can get an interesting result by choosing the U's in the representation Theorem A2.7.1 to be the elements of the chief basis.

$$\Delta\left(\underline{e}',\underline{e}^{2},...,\underline{e}^{n}\right)=\mathcal{E}_{\sigma_{1}\sigma_{2}...\sigma_{n}}\left(\underline{e}'\right)^{\sigma_{1}}\left(\underline{e}^{2}\right)^{\sigma_{2}}...\left(\underline{e}^{n}\right)^{\sigma_{n}}\Delta\left(\underline{e}_{1},\underline{e}_{2},...,\underline{e}_{n}\right).$$

By Theorem A2.5.8 and Daginitim A2.5.3,

$$(\underline{e}')^{\sigma_i} = \underline{e}^{\sigma_i} \cdot \underline{e}' = g^{\sigma_i 1}$$
, etc.

 $\Delta(e',e^2,\dots,e^n) = \varepsilon_{\sigma,\sigma_1\dots\sigma_n} g^{\sigma,1} g^{\sigma,2} \dots g^{\sigma,n} \Delta(e_n,e_2,\dots,e_n).$ $det [g^i \delta]$

The last step is essentially the column expansion definition of the Leterminant of a square matrix. By this and "See, e.g., HOHN's Elementary Matrix Algebra.

Theorem A2.5.10 (vii), we are lead to

Theorem A 2.7.6. Let {e, e, -, en} and {e', e², -, en}

Be a basis and its sual for an oriented n-dimensional

inner product space Vn. Then

(i)
$$\frac{\Delta(\underline{e}',\underline{e}',-\cdot,\underline{e}'')}{\Delta(\underline{e}_1,\underline{e}_2,-\cdot,\underline{e}_n)} = \text{let}[\underline{g}'i] = \frac{1}{\text{let}[\underline{g}_{ij}]};$$

(ii)
$$\Delta(\underline{e}_1,\underline{e}_2,-\cdot,\underline{e}_n)\Delta(\underline{e}',\underline{e}^2,-\cdot,\underline{e}^n)=1$$

(iii)
$$\left[\Delta(\underline{e}_1,\underline{e}_2,...,\underline{e}_n)\right]^2 = \det[g_{ij}], \left[\Delta(\underline{e}_1,\underline{e}_2,...,\underline{e}_n)\right]^2 = \det[g_{ij}].$$

Proof, (i): This part was established by the remarks proceeding the statement of the theorem.

(ii): Let & Exis, Exis, -.., Eens } be an arthonormal basis for Vn. Then by Theorems A2,7.1 and A2.3.14,

$$\Delta(\underline{e}_1,\underline{e}_2,\cdots,\underline{e}_n)\Delta(\underline{e}',\underline{e}^2,\cdots,\underline{e}^n) =$$

$$= \varepsilon_{\sigma_{1}\sigma_{2}\cdots\sigma_{n}} \alpha_{1}^{\sigma_{1}} \alpha_{2}^{\sigma_{2}} \cdots \alpha_{n}^{\sigma_{n}} \Delta(\underline{e}_{(1)}, \underline{e}_{(2)}, \cdots, \underline{e}_{(n)}) \times$$

where $a_i \delta = \underline{e}_{ij}, \underline{e}_i$ and $b^i \delta = \underline{e}_{ij}, \underline{e}^i$. By Axim (03),

This particular proof was shown to the author by Dr. Seyoung Im.

so we are left with

$$\Delta(\underline{e}_1,\underline{e}_2,\cdots,\underline{e}_n)\Delta(\underline{e}',\underline{e}^2,\cdots,\underline{e}^n)=$$

$$=\underbrace{\varepsilon_{\sigma,\sigma_{1}\cdots\sigma_{n}}}_{\text{df}}\underbrace{a_{i}^{\sigma_{1}}a_{z}^{\sigma_{2}}\cdots a_{n}^{\sigma_{n}}}_{\text{lef}}\underbrace{\varepsilon_{\tau_{i}\tau_{2}\cdots\tau_{n}}}_{\text{bf}}\underbrace{b^{l\tau_{i}}b^{2\tau_{2}}\cdots b^{n\tau_{n}}}_{\text{lef}}$$

= Let A Let B

- Let A Let BT

= Let (ABT)

$$= \det \left[\sum_{j=1}^{n} a_{i} \delta b^{k} \delta \right],$$

where have used some obvious matrix notation as well as some standard results about determinants of matrices. I Unraveling this notation, we get

$$\sum_{j=1}^{n} a_{i} \delta_{b} b^{k} \delta_{j} = \sum_{j=1}^{n} (e_{x_{j}} e_{i}) (e_{x_{j}} e_{j}) \quad (\text{substitutum})$$

$$= e_{i} e_{k} \quad (\text{Theorem A 2.3.16})$$

$$= \delta_{i}^{k} \quad (\text{Theorem A 2.5.10 (ii)}).$$

'See, e.g., HOHN's Elementary Matrix Algebra.

Since det [Sik] = 1, we have the desired result.

(iii): These results follow Livertly from countining (i) and (ii). [

The following observation is occasionally wefore.

Theorem A2.7.7, Let { Exis, Exis, exis} be a positive orthonormal basis for V_n . Then for { U_1 , U_2 , ..., U_n } $\subset V_n$,

 $\Delta(\underline{u}_{i},\underline{u}_{z},...,\underline{u}_{n}) = \det[\underline{u}_{i}^{\langle j\rangle}],$ $\underline{u}_{i} = \underline{u}_{i}^{\langle j\rangle}\underline{e}_{zj\rangle},$

Roof, By Theorem A2.7.2 and Definition A2.7.6,

$$\begin{split} \Delta(\mathcal{Q}_{1},\mathcal{Q}_{2},\cdots,\mathcal{Q}_{n}) &= \mathcal{E}_{\mathcal{T}_{1}\mathcal{T}_{2}\cdots\mathcal{T}_{n}} \mathcal{U}_{1}^{\langle \mathcal{T}_{1}\rangle} \mathcal{U}_{2}^{\langle \mathcal{T}_{2}\rangle} \cdots \mathcal{U}_{n}^{\langle \mathcal{T}_{n}\rangle} \\ &= \mathcal{L}et\left[\mathcal{U}_{i}^{\langle \mathcal{S}^{\prime}\rangle}\right] \quad \left(\mathcal{L}_{fini}him\right). \quad \Box \end{split}$$

We are now in a good position to establish the following useful relationship between alternating symbols and Kronecken deltas in the "three-dimensiones" case. Next that despite our method of proof, the result stands by itself independent of the concept of a linear space,

Theorem A 2.7.8. For i, j, k, p, q, r ∈ II3,

(iii)
$$\sum_{i=1}^{3} \sum_{j=1}^{3} \varepsilon_{ijk} \varepsilon_{ijr} = 2 S_{kr}$$

(iv)
$$\sum_{i=1}^{3}\sum_{j=1}^{3}\sum_{k=1}^{3}\epsilon_{ijk}\epsilon_{ijk}=6$$
.

Proof. (i): Let & ELIX, ELZX, ELZX? be a positive orthonormal basis for Vz. Then

$$\Delta\left(\mathcal{L}_{cir},\mathcal{L}_{cjr},\mathcal{L}_{zer}\right) = \mathcal{E}_{ijk} \Delta\left(\mathcal{L}_{cir},\mathcal{L}_{zer},\mathcal{E}_{cer}\right) \left(\mathcal{L}_{cir},\mathcal{L}_{cer},\mathcal{L}_{cer}\right)$$

$$= \mathcal{E}_{ijk} \left(positiveness of \{\mathcal{L}_{cir},\mathcal{L}_{cir},\mathcal{L}_{cir},\mathcal{L}_{cir}\}\right)$$

Thus, by the previous theorem,

$$\mathcal{E}_{ijk} = \det \begin{bmatrix} (e_{zi})^{n} (e_{zi})^{2} (e_{zi})^{3} \\ (e_{zj})^{(i)} (e_{zj})^{(2)} (e_{zj})^{3} \\ (e_{zj})^{(i)} (e_{zk})^{(i)} (e_{zk})^{3} \end{bmatrix}$$

$$= \det \begin{bmatrix} \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{zij} \cdot \mathcal{C}_{zj} & \mathcal{C}_{xij} \cdot \mathcal{C}_{zj} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{zij} \cdot \mathcal{C}_{zj} & \mathcal{C}_{xij} \cdot \mathcal{C}_{zj} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{zij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{zij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} & \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \cdot \mathcal{C}_{xij} \\ \mathcal$$

$$\begin{aligned} \mathcal{E}_{ijk} & \mathcal{E}_{pqr} = \det \begin{bmatrix} \delta_{i1} \delta_{i2} \delta_{i3} \\ \delta_{j1} \delta_{j2} \delta_{j3} \end{bmatrix} \det \begin{bmatrix} \delta_{p1} \delta_{p2} \delta_{p3} \\ \delta_{q1} \delta_{q2} \delta_{q3} \\ \delta_{k1} \delta_{k2} \delta_{k3} \end{bmatrix} \end{aligned}$$

= let
$$\left(\frac{3}{2} S_{im} S_{pm}\right) \left(\frac{3}{2} S_{im} S_{qm}\right) \left(\frac{3}{2} S_{im} S_{rm}\right) \left(\frac{3}{2} S_{im} S_{pm}\right) \left(\frac{3}{2} S_{im}\right) \left(\frac{3}{2} S_{im$$

(Substitution property of Kronecher Lelta)

(ii): Expanding the converteterminant, we have $\begin{aligned} & \text{Eijk Epqr} = \text{Sip}(\text{Sjq Sbr} - \text{Skq Sjr}) - \text{Siq}(\text{Sjp Shr} - \text{Shp Sjr}) \\ & + \text{Sir}(\text{Sjp Skq} - \text{Skp Sjq}); \end{aligned}$

and then

$$\frac{3}{\sum_{i=1}^{3} \epsilon_{ijk}} \epsilon_{iqr} = \frac{3}{i=1} \left[\delta_{ii} \left(\delta_{jq} \delta_{kr} - \delta_{kq} \delta_{jr} \right) - \delta_{iq} \left(\delta_{ji} \delta_{kr} - \delta_{ki} \delta_{jr} \right) \right] \\
+ \delta_{ir} \left(\delta_{ji} \delta_{kq} - \delta_{ki} \delta_{jq} \right) \\
= \sum_{i=1}^{3} \left[\delta_{ii} \left(\delta_{jq} \delta_{kr} - \delta_{kq} \delta_{jr} \right) - \left(\delta_{iq} \delta_{ji} \delta_{kr} + \left(\delta_{iq} \delta_{kr} \right) \delta_{jr} \right) \right] \\
+ \left(\delta_{ir} \delta_{ji} \right) \delta_{kq} - \left(\delta_{ir} \delta_{ki} \right) \delta_{jq} \left[\left(\omega \omega \omega \right) \right]$$

(iii), (iv): Exercise A2.7,5,

Supplementary Reading

BORISENKO and TARAPOV, Vector and Tensor Analysis with
Applications

GREUB, Linear Algebra

MARTIN and MIZEL, Introduction to Linear Algebra

NICKERSON, SPENCER, and STEENROO, Advanced Calculus

A2.8. The Vector Product

Almost certainly you have had some experience with the vector or cross product of two vectors from three-dimensional Endidorn space. Elementary treatments of the vector product often have the difect that somewhere in the development loose or imprecise notions such as "right-hand rule" or "like-handed triads" (usually explained with sketches of the author's extremities) are snuck in to give the direction of the vector product. In this section, we shall use the concept of volume unautation functionals to introduce the vector product in a mathematically sound way-

The vector product of two vectors is strictly a three-vinewind convert, and accordingly for the whole of this section V_3 denotes an oriented three-dimensional inner product space? The associated volume orientation functional is Δ .

We get the cross product from the volume orientation functional as follows. Let $\underline{u}, \underline{x}, \underline{w} \in V_3$. Then for fixed \underline{u} and \underline{x} , $\Delta(\underline{u}, \underline{x}, \underline{w})$ is a linear functional in \underline{w} . More precisely, $\Delta(\underline{u}, \underline{x}, \underline{\cdot})$ is a linear function on V_3 to R. By the representation theorem for linear functionals \overline{u} , \overline{u} a unique vector, which we

G. Theorem AZ. 7. 2 and Definition AZ. 7.4.

² g. Theorem AZ.3.Z,

³ Theorem A2.5.1.

Inste by $\underline{u} \times \underline{v}$ since if can be expected to depend on \underline{u} and \underline{v} , $\exists \Delta(\underline{u},\underline{v},\underline{w}) = (\underline{u} \times \underline{v}) \cdot \underline{w} \quad \forall \ \underline{w} \in V_3$. Hence, we have established

Theorem A2.8.1. To each ordered pain (u, v) of elements from V_3 there corresponds a unique element $u \times v \in V_3 \Rightarrow$

(w, x, w) = (axx)·w ∀w∈1/3.

ux is called the vector product of a and v.

Obviously, the vector product depends on the choice of volume orientation functional. By Theorem A2.7.2, this amounts to changing the vector product by a scalar factor of -1 (1.e., reversing its "sense") if the alternative orientation for V3 is chosen. This is where "right-hand rules", etc. cuter into the elementary geometrical treatment.

Insert material on p. AZ.B.Za

The following result provides component representations for the vector product.

Theorem A2, B, Z, Let $\{\underline{e}_1,\underline{e}_2,\underline{e}_3\}$ and $\{\underline{e}',\underline{e}^2,\underline{e}^3\}$ be a lasis and its Sual for V_3 , Then $\forall \underline{u},\underline{v} \in V_3$,

 $u \times v = \Delta(\mathcal{Q}_1, \mathcal{Q}_2, \mathcal{Q}_3) \varepsilon_{ijk} u j v^k \mathcal{Q}^i$

[&]quot;Often, cross product. Sometimes the notation UNX is used, but we shall save this for the "skew product" of § A3.6.

INSERT for p. AZ.B. 2

Exercise A2,8.0. Prove that in terms of components relative to a positive or thonormal basis

$$(\underbrace{\omega^{\vee} v}) \cdot \underline{w} = \begin{vmatrix} u^{\langle 17} & u^{\langle 27} & u^{\langle 37} \\ v^{\langle 17} & v^{\langle 27} & v^{\langle 27} \end{vmatrix}, \\ w^{\langle 17} & w^{\langle 27} & w^{\langle 37} \end{vmatrix},$$

return to p. A2.3.2

= $\Delta(\xi', \xi^2, \xi^3) \xi^{ijk} (\xi_i \vee_k \xi_i .)$

In particular, in terms of a positive orthonormal basis $\{e_{(1)}, e_{(2)}, e_{(2)}, e_{(3)}\} = \{e^{(1)}, e^{(2)}, e^{(3)}\},$

uxv = Eijk u (i) v (k) e (i).

Proof. By Theorem A2.7.1,

 $\Delta(\underline{\omega},\underline{v},\underline{w}) = \varepsilon_{ijk} \, u^i v \dot{\delta} w^k \, \Delta(\underline{e}_1,\underline{e}_2,\underline{e}_3).$

On defining $x_k = \Delta\left(\underline{e}_1,\underline{e}_2,\underline{e}_3\right) \, \epsilon_{ijk} \, u^i \, v \delta$, we have $\Delta(\underline{u},\underline{v},\underline{w}) = x_k \, w^k$

.. by Theorem A 2.5.10 (iii),

 $\Delta(\underline{w},\underline{v},\underline{w}) = \underline{x} \cdot \underline{w},$

where $x := x_k e^k$. Then by Theorem A2.3.1, $x = u \times v$. To get exact agreement with the assertion of the theorem, we note that by the skew-symmetry of the alternating symbol?

¹ the, & is simply the three-dimensional afterning symbol withen with superscripts for the sake of the summation concertion.

² See Definition A2.7.2.

$$\begin{split} x_k &= \Delta(\underline{e}_1,\underline{e}_2,\underline{e}_3) \, \varepsilon_{ijk} \, u^i v \delta = -\Delta(\underline{e}_1,\underline{e}_2,\underline{e}_3) \, \varepsilon_{ikj} \, u^i v \delta \\ &= + \Delta(\underline{e}_1,\underline{e}_2,\underline{e}_3) \, \varepsilon_{kij} \, u^i v \delta \;, \end{split}$$

To get the formula for the contravariant components of UXV start with

 $\Delta(\underline{u},\underline{v},\underline{w}) = \varepsilon^{ijk} u_i v_j w_k \Delta(\underline{e}',\underline{e}^2,\underline{e}^3)$

and follow the same steps. I

We can use the above component representation of the vector product to get a geometrical interpretation which is the austomany elementary befinition. First we turn to some geometrical preliminaries.

Definition A2.8.1. Let u and x be any two un parallel elements of 73. Then the plane of u and y is Lsp Eu, x ?. A unit normal to the plane of u and x is an element n e V3 with the properties that

n·w=0 \welsp{u,x} and In 1=1,

The following result is intuitively evident; but, as is usually the case in such matters, its proof is rather

In particular, reither w nor & can be the zero element of 1/3; g. Definition AZ-3.6 and Theorem AZ.Z.3.

involved.

Theorem AZ. B. 3. Let u and v be two nonparallel elements of V3. Then the plane of is and x

(i) is a two-dimensional linear subspace of V3

(ii) has exactly two unit normals, with one the regative of the other. [Lsp {u, x} is a linear subspace by Theorem A2.1.12.

Proof, (i) : Since wand & are nonpurallel, the set { w, x } is linearly independent by Sequistion. Also by construction, {u, x} spans Lsp {u, x}. Hence, {u, x} is a losis for Lsp {u, x}, and the dimension of Lsp {u, x} is Z.)

(ii): For notational convenience, set u= E, v= Ez. By Theorem A2, 2,9, {e, e, } can be extended to a casis for 1/3; i.e., ∃ e, ∈ V3 → {e, e, e, e, } is a lasis for 1/2.

Now we want to find an n e V3 >

(necessary and) It is sufficient to have $\eta \cdot \varepsilon_1 = \eta \cdot \varepsilon_2 = 0$, Then by Theorem A2.5.8 and Definition A2.5.2,

See Definition A2.2.2. 2 Prove this as Exercise A2.8.1.

$$\left.\begin{array}{l}
\eta_{1} = \underline{e}_{1} \cdot \underline{n} = 0 \\
\eta_{2} = \underline{e}_{2} \cdot \underline{n} = 0
\end{array}\right\} \implies \underline{n} = \underline{n}_{3} \underline{e}_{3},$$

where $\{e',e^2,e^3\}$ is the dual Gasis of $\{e_1,e_2,e_3\}$. The condition $|n|=1 \Rightarrow |n_3|=(|e^3|)^{-1}$. Thus,

$$\mathcal{D} = \pm \frac{1}{|\mathcal{E}^3|} \mathcal{E}^3.$$

At this point, we have two n's, and one is the regative of the other, However, there may well be other pairs of unit mounts because of the nonuniqueness of the element \mathcal{E}_3 chosen to extend $\{\mathcal{E}_1,\mathcal{E}_2\}$ to a tosis. To investigate this, suppose that the basis is completed by appending a different element, say \mathcal{E}_3 . The dual basis of $\{\mathcal{E}_1,\mathcal{E}_2,\mathcal{E}_3\}$ is denoted by $\{\mathcal{E}_1,\mathcal{E}_2,\mathcal{E}_3\}$. In terms of the original dual basis $\{\mathcal{E}_1,\mathcal{E}_2,\mathcal{E}_3\}$ is denoted by $\{\mathcal{E}_1,\mathcal{E}_2,\mathcal{E}_3\}$. In terms of the original dual basis $\{\mathcal{E}_1,\mathcal{E}_2,\mathcal{E}_3\}$, the element \mathcal{E}_3 has the representation

$$\overset{*}{\mathcal{E}}^{3} = \alpha_{1} \overset{e}{\mathcal{E}}^{1} + \alpha_{2} \overset{e}{\mathcal{E}}^{2} + \alpha_{3} \overset{e}{\mathcal{E}}^{3}.$$

By Theorems A 2.5, 8 and A 7.5, 4,

$$a_1 = e_1 \cdot e^3 = 0$$
, $a_2 = e_2 \cdot e^3 = 0$;

and i.

$$\overset{*}{\xi}^3 = a_3 \, \overset{*}{\xi}^3 \, ,$$

thus,

$$\frac{1}{|\underline{\mathring{e}}^3|} \overset{\text{def}}{\overset{\text{def}}{=}} = \frac{\alpha_3}{|\alpha_3|} \frac{1}{|\underline{e}^3|} \overset{\text{def}}{\overset{\text{def}}{=}} = \pm \frac{1}{|\underline{e}^3|} \overset{\text{def}}{\overset{\text{def}}{=}} \overset{\text{def}$$

and we see that no matter how { ?, e, ? is extended to a basis, the resulting pair of n's is always the same. \square

Now we are in a position to state and prove the following theorem which corresponds to the elementary geometrical definition of the vector product.

Theorem A2.8.4. Let u and v be arbitrary clements of V_3 and let Θ be the angle between them. Then

 $\frac{(i)}{\omega \times x} = |\omega| |x| \sin \theta \qquad)^{t}$ $\frac{\sin x}{\omega} = |\omega \times x| \sin \theta \qquad)^{t}$

Proof. (i): This assertion follows casily from Theorem A2, B.Z with the aid of the identity

^{&#}x27;g course, if either U=Q or Y=Q, then the angle Θ is not defined; but in this case Theorem A2.8.2 \Rightarrow $U\times Y=Q$.

If u and v, are parallel, they so not span a plane; but in this case again Theorem A2.8.2 $\Rightarrow u \times v = Q$. Work out the setails of the argument as Exercise A2.8.2. The result also follows from $(u \times v) \cdot w = \Delta(u, v, w) = 0 \ \forall \ w \ \text{since} \ \{u, v\} \ \text{is linearly supercent}.$

 $\sum_{i=1}^{3} \mathcal{E}_{ijk} \mathcal{E}_{ilm} = \mathcal{S}_{je} \mathcal{S}_{km} - \mathcal{S}_{jm} \mathcal{S}_{ke} \quad (\textit{Theorem A2.7,8 (ii)}).$

The Setails are left as Exercise A2.8.3. Alternatively, one can get this result by combining the representation of Theorem A2.8.2 with the special choice of basis used in Exercise A2.7.3, Work out such a perof as Exercise A2.8.4.

 $\Delta(\mathcal{U}, \mathcal{X}, \mathcal{D}_2) = \Delta(\mathcal{U}, \mathcal{X}, -\mathcal{D}_1) \qquad (\text{substitution})$ $= -\Delta(\mathcal{U}, \mathcal{X}, \mathcal{D}_1) \; ; \qquad (\text{multiluneauty})$

and only one of $\Delta(\underline{u},\underline{v},\underline{n},)$ and $\Delta(\underline{u},\underline{v},\underline{n}_2)$ is positive.

By Theorem AZ.2.8, we can take $\{u, v, p\}$ to be a basis for V_3 . For convenience, we write

$$e_{j}=\alpha \Rightarrow \alpha'=1, \alpha^{2}=\alpha^{3}=0;$$

$$\mathcal{E}_{3} := n \Rightarrow n' = n^{2} = 0, \quad n^{3} = 1.$$

Substitution of this information into the component foundary Theorem AZ.B.Z yields

$$\mathcal{L} \times \mathcal{L} = \Delta(\mathcal{L}, \mathcal{L}, \mathcal{L}) e^{3}$$
.

By Theorem A2.5.5, & is uniquely determined by the conditions

$$e^3 \cdot e_1 = e^3 \cdot u = 0$$

$$e^{3} \cdot e_{2} = e^{3} \cdot v = 0$$

$$e^3 \cdot e_3 = e^3 \cdot n = 1$$
.

Clearly, $e^3 = n$ satisfies these conditions; and by the uniqueness, nothing else will. Hence, we have

$$\widetilde{u} \times \widetilde{v} = \Delta(\widetilde{u}, v, \widetilde{u}) \widetilde{u}$$
.

$$| \mathcal{U}_{\times} \times \mathcal{V}_{\times} | = | \Delta(\mathcal{U}_{\times}, \mathcal{V}_{\times}, n) | | \mathcal{U}_{\times} | \quad \text{(Theorem A 2.3.6 (c))}$$

$$= \Delta(\mathcal{U}_{\times}, \mathcal{V}_{\times}, n) \quad \left(\Delta(\mathcal{U}_{\times}, \mathcal{V}_{\times}, n) > 0 \text{ and } | \mathcal{U}_{\times} | = 1 \right),$$
and :.
$$\mathcal{U}_{\times} \times \mathcal{V}_{\times} = | \mathcal{U}_{\times} \times \mathcal{V}_{\times} | \mathcal{U}_{\times},$$

(iii): Established above, [

Exercise A2.8.6. Use Theorem A2.8.4 to show that $|\Delta(i\underline{v},\underline{v},\underline{w})| = |(u\underline{x}\underline{v})\cdot\underline{w}|$ is the volume (in the usual geometrical sense) of the parallelepiped generated by $\underline{u},\underline{x}$, and \underline{w} , Do not hositate to draw pichures.

The next theorem gathers some of the rules for manipulating the vector product. Each rule is accessible from Theorem A2.8.1, our abstract definition of the vector product in terms of the volume orientation functional, or Theorem A2.8,2, the representation of the vector product in terms of components, or Theorem A2.8.4, the geometrical interpretation of the vector product, but most of them are established most readily via the component representation.

Theorem A2. B. 5. The vector product has the following properties:

(i) Shew-Symmetry.
$$\forall u, v \in V_3$$

 $v \times u = -(u \times v)$;

- (ii) Homogeneity. $\forall u, x \in V_3 \text{ and } \forall \alpha, \beta \in \mathbb{R}$ $(\alpha u) \times (\beta x) = (\alpha \beta) (u \times x);$
- (iii) Distributivity w.r.t. addition. $\forall u, x, w \in V_3$ $u \times (v + w) = u \times v + u \times w$

(x+x)xx = xxx +xxx ;

(iv) $\forall \underline{\omega}, \underline{v}, \underline{w} \in V_3$

(u==) × w - (u·w) × - (v.w) w ;)'

(v) <u>For</u> u, y ∈ V₃

Ux V = Q iff u and x are parallel;

(Vi) For u ∈ V3

Proof. Exercises AZ. 8.7-12. [

We and this section with an inkesting application of the

Thus, in general, $(u \times v) \times w \neq u \times (v \times w)$; i.e., the vector product is monassociative.

the component representation formula of Theorem A2.B.2. To wit,

$$\mathcal{E}_{2} \times \mathcal{E}_{3} = \Delta(\mathcal{E}_{1}, \mathcal{E}_{2}, \mathcal{E}_{3}) \, \mathcal{E}_{ijk} \left(\mathcal{E}_{2}\right)^{k} \left(\mathcal{E}_{3}\right)^{k} \mathcal{E}^{i} = \Delta(\mathcal{E}_{1}, \mathcal{E}_{2}, \mathcal{E}_{3}) \mathcal{E}^{i},$$

$$\mathcal{E}_{i23}^{i} \quad \mathcal{E}_{i23}^{k}$$

which provides a formula for calculating & . Formally, we have established

Theorem A2.8.6. Let {e, e, e, e, } and {e, e, e, e} be a basis and its dual for V3. Then

$$\underline{e}' = \frac{\underline{e}_z \times \underline{e}_3}{\Delta(\underline{e}_1, \underline{e}_z, \underline{e}_3)} , \underline{e}^2 = \frac{\underline{e}_3 \times \underline{e}_1}{\Delta(\underline{e}_1, \underline{e}_z, \underline{e}_3)} , \underline{e}^3 = \frac{\underline{e}_1 \times \underline{e}_2}{\Delta(\underline{e}_1, \underline{e}_z, \underline{e}_3)} .$$

Supplementary Reading Same as for \$ 42.7